INTRODUCTION TO VLSI SYSTEMS

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PREFACE:

As a result of improvements in fabrication technology, Large Scale Integrated (LSI) electronic circuitry has become so dense that a single silicon LSI chip may contain tens of thousands of transistors. Many LSI chips, such as microprocessors, now consist of multiple, complex subsystems, and thus are really integrated systems rather than integrated circuits.

What we have seen so far is only the beginning. Achievable circuit density now approximately doubles with each passing year. How long can this continue, and how small can the transistor be made? From the physics we find that the linear dimensions of transistors can be reduced to less than 1/10 of those in current integrated systems before they cease to function as the sort of switching elements from which we can easily build digital systems. It will eventually be possible to fabricate chips with hundreds of times as many components as today's. The transistors in such very large scale integrated (VLSI) systems will ultimately have linear dimensions smaller than the wavelength of visible light. Now, non-optical, high-resolution lithographic techniques are now being developed by many firms to enable fabrication of such circuitry.

The emerging VLSI presents a challenge not only to those involved in development of appropriate fabrication technology, but also to computer scientists and computer architects. The ways in which digital systems are structured, the procedures used to design them, the tradeoffs between hardware and software, and the methodology and metrics of analysis of algorithms will all be greatly affected by the maturation of electronic technology towards its maximum density. We believe this will be an important area of activity for computer science on through the 1980's.

Until recently, the design of integrated electronic circuitry has been largely the province of circuit and logic designers working within semiconductor firms. Computer system architects have traditionally built systems from standard integrated circuits designed and manufactured by these firms, but haven't often participated in the specification and design of these integrated circuits. Electrical Engineering and Computer Science (EE/CS) curricula have reflected this tradition, with courses in device physics and integrated circuit design aimed at and generally taken by different students than those interested in digital system architecture and computer science.

This text is written to fill a current gap in the literature, and provide students of computer science and electrical engineering with an introduction to integrated system architecture and design. Combined with individual study in related research areas and participation in actual system design projects, the text could serve as a basis for a graduate course sequence in integrated systems. Portions could be used for an undergraduate text on the subject, or to augment a graduate course on computer architecture. It could also be used to extend, in the system direction, a classical electrical engineering course in integrated circuits. We assume the reader's background contains the equivalent of an introductory course sequence in computer science, and introductory courses in electronic circuits and digital design.
Up till now there have been major obstacles in the path of those attempting to gain an overall understanding of integrated systems. Integrated electronics, developing in a heatedly competitive business environment, has proliferated into a large array of different device technologies, circuit design families, logic design techniques, maskmaking and wafer fabrication techniques, etc. The technologies have sprung up from the grass roots of "Silicon Valley" in California. Most participants in the industry have of necessity concentrated on rather narrow specialties. Texts on the subject have tended to give detailed accounts of some very narrow horizontal segment of the overall subject, such as device physics or circuit design.

We have chosen instead to provide the minimum of basic information about devices, circuits, fabrication technology, logic design techniques, and system architecture, which is sufficient to enable the reader to span fully the entire range of abstractions, from the underlying physics to complete VLSI digital computer systems. A rather small set of carefully selected key concepts is all that is necessary for this purpose. We believe that only by carrying along the least amount of unnecessary mental baggage at each step in such a study, will the student emerge with a good overall understanding of the subject. Once this range of abstractions is spanned, the sequence of concepts can then be mapped into the reader's own space of application and technology.

Another major obstacle has been the high rate of change of integrated electronics. The uninitiated could easily get the feeling that much energy could be invested in learning material which becomes obsolete as rapidly as it is assimilated. The major driving mechanism in all this change is the continual improvement in fabrication technology. This evolutionary process results in the feasibility of manufacturing smaller and smaller devices as time passes. By including the effects of scaling down device dimensions as an essential ingredient of all topics in this text, many of the important changes of the architectural parameters of the technology are predicted, expected, and indeed hoped for.

The key concepts of this text are illustrated by way of specific examples. In any given technology, form follows function in a particular way. The skill of mapping function into form, when once acquired, can be readily applied to any technology. Because of its density, speed, topological properties, and general availability of wafer fabrication, nMOS has been chosen as the technology in which examples are implemented.

An atmosphere of excitement and anticipation pervades this field. Workers from many backgrounds, computer scientists, electrical engineers, and physicists, are collaborating on a common problem area which has not yet become classical. The territory is vast, and largely unexplored. The rewards are great for those who simply press forward.

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BACKGROUND

Prior to the commercial publication of this textbook, this limited printing is being distributed to a selected group of universities as course notes for graduate courses on integrated systems, for the purpose of obtaining critical reviews. Copies are also being distributed to the industrial participants in the Caltech Silicon Structures Project, and to selected individuals in universities and in industry for their review. The authors welcome any and all comments and suggestions from readers. We are especially interested in hearing of the experiences of those teaching from the text. Notifications and corrections of errors, ideas for improvements in the tutorial techniques used, and suggestions of instructive problems for each chapter would be greatly appreciated.

While the material in this text is presented in a particular order, it need not be read in that order. Each chapter presents material from a distinct level in the hierarchy of disciplines involved in integrated systems. The material falls into four major groupings: chapters 1 and 2 provide the basics of devices and fabrication, chapters 3 and 4 give the basics of design and implementation, chapters 5 and 6 present an example of LSI system design, and chapters 7, 8, and 9 discuss topics of current interest in integrated systems research. We recommend that the reader start in the chapter where he or she is most knowledgeable, and read until information is required from an adjacent area described in some other chapter. By using this algorithm and consulting the reading references where necessary, the reader can gradually work through the primary material of all chapters. Although much of the material in this text is previously unpublished, it contains only fundamental concepts. However, these concepts cover a wide range of disciplines, and are easily visualized only after the overall context becomes clear.

This text has its origins in a series of courses in integrated circuit design given by Carver Mead at Caltech, beginning in 1970. Starting in 1971, students in these courses designed and debugged their own integrated circuits. The students undertook increasingly complex system designs, using only rather simple implementation aids. The structured design methodology presented in this text evolved within this milieu. A separate Computer Science department was formed at Caltech in 1976, with integrated systems as a focus. An early association was formed with systems architects in industry. Interaction of Caltech students and faculty with industrial researchers resulted in the idea of a joint Caltech-industry cooperative program. The initial industrial participants in this program, known as the Caltech Silicon Structures Project, are Intel, DEC, IBM, Hewlett-Packard, Xerox, and Burroughs.

Work on this text began in August 1977. The first three chapters were used as course notes during the fall of 1977, in courses given by Carver Mead at Caltech and by Carlo Sequin at U. C. Berkeley. The first five chapters were used during the spring of 1978 in courses given by Ivan Sutherland at Caltech, by Robert Sproull at Carnegie-Mellon University, and by Fred Rosenberger at Washington University, St. Louis. This printing is being used again during the fall of 1978 in the courses at Caltech and U.C. Berkeley, and in new courses by Lynn Conway, while visiting at M.I.T., by Kent Smith at University of Utah, and by Bob Bower at U.C.L.A.
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Chapter 1: MOS Devices and Circuits

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Sections:

The MOS Transistor • • • The Basic Inverter • • • Inverter Delay • • • Parasitic Effects • • •
Driving Large Capacitive Loads • • • Space vs Time • • • Basic NAND and NOR Logic
Circuits • • • Super Buffers • • • A Closer Look at the Electrical Parameters • • • Depletion Mode
vs Enhancement Mode Pullups • • • Delays in Another Form of Logic Circuitry • • •
Pullup/Pulldown Ratios for Inverting Logic Coupled by Pass Transistors • • • Transit Times and
Clock Periods • • • Properties of Cross Coupled Circuits • • • A Fluid Model for Visualizing MOS
Transistor Behavior • • • Effects of Scaling Down the Dimensions of MOS Circuits and Systems

In this chapter we begin with a discussion of the basic properties of the n-channel, metal-oxide-semiconductor (MOS), field effect transistor (FET). We then describe and analyze a number of
circuits composed of interconnected MOS field effect transistors. The circuits described are
typical of those we will commonly use in the design of integrated systems. The analysis, though
highly condensed, is conceptually correct and provides a basis for the solution of most system
problems typically encountered.

Integrated systems in MOS technology contain three levels of conducting material separated by
intervening layers of insulating material. Proceeding from top to bottom, these levels are termed
the metal, the polysilicon, and the diffusion levels respectively. Patterns for paths on these three
levels, and the locations of contact cuts through the insulating material to connect certain points
between levels, are transferred into the levels during the fabrication process from masks similar to
photographic negatives. The details of the fabrication process will be discussed in chapter 2.

In the absence of contact cuts through the insulating material, paths on the metal level may cross
over paths on the polysilicon or diffusion levels with no significant functional effect. However,
wherever a path on the polysilicon level crosses a path on the diffusion level, a transistor is
created. Such a transistor has the characteristics of a simple switch, with a voltage on the
polysilicon level path controlling the flow of current in the diffusion level path. Circuits
composed of such transistors, interconnected by patterned paths on the three levels, form our
basic building blocks. With these basic circuits, we will architect integrated systems, to be
fabricated on the surface of monolithic crystalline chips of silicon.
The MOS Transistor

An MOS transistor will be produced on the integrated system chip wherever a polysilicon path crosses a diffusion path, as shown in figure 1a. The electrical symbol used to represent the MOS transistor in our circuit diagrams is shown in figure 1b, along with symbols and polarities of certain voltages of interest. Note that the source and drain terminals of the device are physically symmetrical. For the n-channel MOSFETs, these terminal labels are assigned such that $V_{ds}$ is normally positive. A more detailed view of the rectangular region called the gate, where the polysilicon (poly) crosses the diffusion, is given in figure 1c. During fabrication the diffusion paths are formed after the poly paths are formed, as explained more fully in chapter 2. The poly gate, and the thin layer of oxide beneath it, mask the region under the gate during diffusion. Therefore, no diffusion path forms under the gate, and there is no direct connection on the diffusion level between the source and drain terminals of the transistor. Notice in this discussion that metal, poly, and diffusion paths all conduct electricity well enough to be considered "wires" until further notice.

In the absence of any charge on the gate, the drain to source path through the transistor is like an open switch. The gate, separated from the substrate by the layer of thin oxide, forms a capacitor. If sufficient positive charge is placed on the gate so that $V_{gs}$ exceeds a threshold voltage $V_{th}$, electrons will be attracted to the region under the gate to form a conducting path between drain and source. Most of the transistors we will use in our systems have threshold voltages greater than zero. These are called enhancement mode MOSFETs, and their threshold voltage typically equals $\sim 0.2(VDD)$, where VDD is the positive supply voltage for the particular technology.

The basic operation performed by the MOS transistor is to use charge on its gate to control the movement of negative charge between source and drain through the channel under the gate. The current from source to drain equals the charge induced in the channel divided by the transit time or average time required for an electron to move from source to drain. The transit time itself is the distance the electron has to move divided by its average velocity. In semiconductors under normal conditions, the velocity is proportional to the electric field driving the electrons. The relationship between drain to source current $I_{ds}$, drain to source voltage $V_{ds}$, and gate to source voltage $V_{gs}$ is sketched in figure 1d. For small $V_{ds}$, the transit time $\tau$ is given by equation 1.
Fig. 1a. MOS Transistor (top view)

Fig. 1b. MOS Transistor Symbol (subscripts in + to - direction sequence)

Fig. 1c. MOSFET Gate Dimensions

Fig. 1d. Current vs Voltage
Transit time: \[ \tau = \frac{L}{\text{velocity}} = \frac{L}{\mu E} = \frac{L^2}{|\mu| V_{ds}} \] [eq.1]

The proportionality constant \( \mu \) is called the mobility of the charge carriers, in this case electrons, under the influence of an electric field in the conducting material of the channel region. It is a velocity per unit electric field (cm²/volt-sec). We shall see that the transit time is the fundamental time unit of the entire integrated system.

The amount of negative charge in transit is just the gate capacitance times the voltage on the gate in excess of the threshold voltage. The capacitance of two parallel conductors of area \( A \), separated by insulating material of thickness \( D \), equals \( \varepsilon A/D \). The proportionality constant \( \varepsilon \) is called the permittivity of the insulating material, and has a simple interpretation. It is the capacitance of parallel conductors of area \( A = 1 \text{ cm}^2 \), separated by a thickness \( D = 1 \text{ cm} \) of the insulator material, and is in the units farad/cm. Therefore, the gate capacitance equals \( \varepsilon WL/D \). Thus the charge in transit is given by eq. 2, and the current is given by eq. 3.

Charge in transit: \[ Q = -C_g(V_{gs} - V_{th}) = -\frac{\varepsilon WL}{D}(V_{gs} - V_{th}) \] [eq.2]

Current: \[ I_{ds} = -I_{sd} = \frac{\text{charge in transit}}{\text{transit time}} = \frac{\mu \varepsilon W}{LD}(V_{gs} - V_{th})(V_{ds}) \] [eq.3]

Note that for small \( V_{ds} \), the drain current is proportional to the source-drain voltage and also to the gate voltage above threshold. Any device with a current through it proportional to the voltage across it may be viewed as a resistor, and in the case of an MOS device with low drain to source voltage, the resistance is controlled by the gate voltage as given in eq. 3a.

\[ \frac{V_{ds}}{I_{ds}} = R = \frac{L^2}{|\mu| C_g(V_{gs} - V_{th})} \] [eq.3a]

In both equations 2 and 3a, \( C_g \) is the gate to channel capacitance of the turned on transistor. In the simple case where this transistor is driving the gate of another one identical to it, the time response of the system will be an exponential with a time constant \( RC_g \), given in equation 4. This time constant is identical to the transit time \( \tau \) given in equation 1.
\[ RC_g = \frac{L^2}{\mu(V_{gs} - V_{th})} = \tau \]  

Although the above equations are greatly simplified, they provide sufficient information to make many design decisions which we will face, and also give us insight into the scaling of devices to smaller sizes. In particular, the transit time \( \tau \) can be viewed as the basic time unit of any system we shall build in the integrated technology. In almost all situations, the fastest operation which we can perform is to transfer a signal from the gate of one MOS transistor onto the gate of another. The transit time is the minimum time in which a charge placed on the gate of one transistor results in the transfer of a similar charge through that transistor's channel onto the gate of a subsequent transistor. For example, to transfer a charge from one transistor onto two transistors identical to it requires a minimum of two transit times. Thus, the transit time of the basic transistor in an integrated system can be viewed as the unit of time in which all other times in the system are scaled. Although it is a somewhat optimistic approximation, we will use \( \tau \) as the primary time metric in calculating the delay through elementary inverting logic stages. More accurate predictions of circuit behavior can be produced using any one of a number of available circuit simulation programs.\(^5,6\)

As \( V_{ds} \) is increased, not all of the drain to source voltage is available for reducing the transit time. Drain voltage in excess of one threshold below the gate voltage creates a short region of high electric field adjacent to the drain which the carriers cross very quickly. The electric field in the major portion of the channel from the source up to this region is proportional to \( V_{gs} - V_{th} \), as shown in figure 1e. For \( V_{ds} > (V_{gs} - V_{th}) \), the drain current becomes independent of \( V_{ds} \). Further increases in \( V_{ds} \) neither increase \( I_{ds} \) nor decrease the transit time. This range of \( V_{ds} \) values is known as saturation.

In saturation:

\[ I_{ds} = \frac{\mu F W (V_{gs} - V_{th})^2}{2LD} \]  

[eq.5]

With the exception of the factor of 2 in the denominator, this equation is similar to equation 3, with the \( V_{ds} \) factor in that equation replaced by its maximum effective value, \( V_{gs} - V_{th} \). The factor of 2 in equation 5 arises from the non-uniformity of the electric field in the channel region when in saturation.\(^1,4\)
The Basic Inverter

The first logic circuit we will describe is the basic digital inverter. Analysis of this circuit is then extended to analysis of basic NAND and NOR logic gates. The inverter's logic function is to produce an output which is the complement of its input. When describing the logic function of circuits in integrated systems, we assign the value logic-1 to voltages equaling or exceeding some defined logic threshold voltage, and logic-0 to voltages less than this threshold voltage.

Were there an efficient way to implement resistors in the MOS technology, we could build a basic digital inverter circuit using the configuration of figure 2a. Here, if the inverter input voltage \( V_{\text{in}} \) is less than the transistor threshold voltage \( V_{\text{th}} \), then the transistor is switched off, and \( V_{\text{out}} \) is "pulled-up" to the positive supply voltage VDD. In this case the output is the complement of the input. If \( V_{\text{in}} \) is greater than \( V_{\text{th}} \), the transistor is switched on and current flows from the VDD supply through the resistor R to GND. If R were sufficiently large, \( V_{\text{out}} \) could be "pulled-down" well below \( V_{\text{th}} \), thus again complementing the input. However, the resistance per unit length of minimum width lines of various available conducting elements is far less than the effective resistance of the switched on MOSFET. Implementing a sufficiently large inverter pullup using resistive lines would require a very large area compared to that occupied by the transistor itself.

To circumvent this problem a depletion mode MOSFET is used as a pullup for the basic inverter circuit, symbolized and configured as shown in figure 2b. In contrast to the usual enhancement mode transistor, the depletion mode transistor has a threshold voltage, \( V_{\text{dep}} \), that is less than zero. During fabrication, one of the masks is used to select any desired subset of transistors in the integrated system for processing as depletion mode transistors. For a depletion mode transistor to turn off, it requires a voltage on its gate relative to its source that is more negative than \( V_{\text{dep}} \). But the depletion mode pullup transistor's gate is connected to its source, and thus it is always turned on. Hence, when the enhancement mode transistor is turned off, for example by connecting zero voltage to its gate, the output of the inverter will be equal to VDD. We will find that for reasonable ratios of the gate geometries of the two transistors, input voltages above a defined logic threshold voltage, \( V_{\text{inv}} \), will produce output voltages below that logic threshold voltage, and vice versa.

The top view of the layout of an inverter on the silicon surface is sketched in figure 2c. It
consists of two polysilicon regions overhanging a path in the diffusion level which runs between VDD and GND. This arrangement forms the two MOS transistors of the inverter. The inverter input A is connected to the poly forming the gate of the lower of the two transistors. The pullup is formed by connecting the gate of the upper transistor to its source. The fabrication details of such connections are described in chapter 2. The output of the inverter is shown emerging on the diffusion level, from between the drain of the pulldown and the source of the pullup. The pullup is a depletion mode transistor, and it is usually several times longer than the pulldown, to achieve the proper inverter logic threshold.

Figures 3a and 3b show the characteristics of a typical pair of MOS transistors used to implement an inverter. The relative locations of the saturation regions of the pullup and pulldown differ in these characteristics, due to the difference in their threshold voltages.

We can use a graphical construct to determine the actual transfer characteristic, $V_{\text{out}}$ vs $V_{\text{in}}$, of the inverter circuit. From figure 2b we see that the $V_{ds(\text{enh})}$ of the enhancement mode transistor equals VDD minus $V_{ds(\text{dep})}$ of the depletion mode transistor. Also, $V_{ds(\text{enh})}$ equals $V_{\text{out}}$. In a steady state and with no current drawn from the output, the $I_{ds}$ of the two transistors are equal. Since the pullup has its gate connected to its source, only one of its characteristic curves is relevant, namely the one for $V_{gs(\text{dep})} = 0$. Taking these facts into account, we begin the graphical solution (fig. 3c) by superimposing plots of $I_{ds(\text{enh})}$ vs $V_{ds(\text{enh})}$, and the one plot of $I_{ds(\text{dep})}$ vs [VDD - $V_{ds(\text{dep})}$]. Since the currents in both transistors must be equal, the intersections of these sets of curves yields $V_{ds(\text{enh})} = V_{\text{out}}$ versus $V_{gs(\text{enh})} = V_{\text{in}}$. The resulting transfer characteristic is plotted in figure 3d.

Studying figures 3c and 3d, consider the effect of starting with $V_{\text{in}} = 0$ and then gradually increasing $V_{\text{in}}$ towards VDD. While the input voltage is below the threshold of the pulldown transistor, no current flows in that transistor, the output voltage is constant at VDD, and the drain to source voltage across the pullup transistor is equal to zero. When $V_{\text{in}}$ is first increased above the enhancement mode threshold, current begins to flow in the pulldown transistor. The output voltage decreases slowly as the input voltage is first increased above $V_{th}$. Subsequent increases in the input voltage rapidly lower the pulldown's drain to source voltage, until the point is reached where the pulldown leaves its saturation region and becomes resistive. Then as $V_{\text{in}}$ continues to increase, the output voltage asymptotically approaches zero. The input voltage at which $V_{\text{in}}$
Fig. 1e. Voltage Profile Across Channel

Fig. 2a. An Inverter

Fig. 2b. The Basic Inverter Circuit Diagram, Logic Symbol, Logic Function

Fig. 2c. Basic Inverter Layout
Fig. 3a. Inverter Pullup Characteristics

Fig. 3b. Inverter Pulldown Characteristics

Fig. 3c. $I_{ds(enh)}$ vs $V_{ds(enh)}$, and $I_{ds(dep)}$ vs $[VDD - V_{ds(dep)}]$.

Fig. 3d. $V_{out}$ vs $V_{in}$ for the Basic Inverter.

$V_{out}$

$V_{DD}$

$V_{in}$

$V_{inv}$

$V_{DD}$

decreasing $L_{pu}/W_{pu} = Z_{pu}$

$L_{pd}/W_{pd} = Z_{pd}$
equals $V_{\text{out}}$ is known as the logic threshold voltage $V_{\text{inv}}$. Figure 3d also shows the effect of changes in the transistor length to width ratios on the transfer characteristics and on the logic threshold voltage. The resistive impedance of the MOS transistor is proportional to the length to width ratio $Z$ of its gate region. Using the subscript $\text{pu}$ for the pullup transistor and $\text{pd}$ for the pulldown transistor: If $Z_{\text{pu}} = L_{\text{pu}}/W_{\text{pu}}$ is increased relative to $Z_{\text{pd}} = L_{\text{pd}}/W_{\text{pd}}$, then $V_{\text{inv}}$ decreases, and vice-versa. The gain, or negative slope of the transfer characteristic near $V_{\text{inv}}$ increases as $Z_{\text{pu}}/Z_{\text{pd}}$ increases. The gain $G$ must be substantially greater than unity for digital circuits to function properly.

**Inverter Logic Threshold Voltage**

The most fundamental property of the basic inverter circuit is its logic threshold voltage, $V_{\text{inv}}$. The logic threshold here is not the same as $V_{\text{th}}$ of the enhancement mode transistor, but is that voltage on the input of the enhancement mode transistor which causes an equal output voltage. If $V_{\text{in}}$ is increased above this logic threshold, $V_{\text{out}}$ falls below it, and if $V_{\text{in}}$ is decreased below $V_{\text{inv}}$, $V_{\text{out}}$ rises above it. The following simple analysis assumes that both pullup and pulldown are in saturation, so that equation 3 applies. Usually the pullup is not quite in saturation, but the following is still nearly correct. $V_{\text{inv}}$ is approximately that input voltage which would cause saturation current through the pulldown transistor to be equal to saturation current through the pullup transistor. Referring to eq.5, we find the condition for equality of the two currents given in eq.6.

Currents equal when:  
\[
\frac{W_{\text{pd}}}{L_{\text{pd}}} (V_{\text{inv}} - V_{\text{th}})^2 = \frac{W_{\text{pu}}}{L_{\text{pu}}} (-V_{\text{dep}})^2, \quad [\text{eq.6}]
\]

or thus when:  
\[
V_{\text{inv}} = V_{\text{th}} - V_{\text{dep}}/[(Z_{\text{pu}}/Z_{\text{pd}})]^{1/2} \quad [\text{eq.6a}]
\]

Here we note that the current through the depletion mode transistor is dependent only on its geometry and threshold voltage $V_{\text{dep}}$, since its $V_{\text{gs}} = 0$. Note that $V_{\text{inv}}$ is dependent upon the thresholds of both the enhancement and depletion mode transistors, and also the square root of the ratio of the $Z = L/W$ of the enhancement mode transistor to that of the depletion mode transistor.
Tradeoffs are possible between these threshold voltages and the areas and current driving capability of transistors in the system's inverters. To maximize \((V_{gs} - V_{th})\) and increase the pulldowns' current driving capability for a given area, \(V_{th}\) should be as low as possible. However, if \(V_{th}\) is too low, inverter outputs won't be driveable below \(V_{th}\), and inverters won't be able to turn off transistors used as simple switches. The original choice of \(V_{th} \sim 0.2VDD\) is a reasonable compromise here.

Similarly, to maximize the current driving capability of pullups of given area, we might set the system's \(V_{dep}\) as far negative as possible. However, eq. 6a shows that for chosen \(V_{inv}\) and \(V_{th}\), decreasing \(V_{dep}\) requires an increase in \(L_{pu}/W_{pu}\) typically leading to an increase in pullup area. The compromise made in this case is usually as follows. The negative threshold of depletion mode transistors is set during fabrication such that with gate tied to source, they turn on approximately as strongly as would an enhancement mode transistor with VDD connected to its gate and its source grounded. In other words, depletion mode and enhancement mode transistors of equal gate dimensions would have equal drain to source currents under those conditions. Applying eq.6 in those conditions we find that:

\[
(-V_{dep})^2 \sim (VDD - V_{th})^2.
\]

Therefore, \(-V_{dep} \sim (VDD - V_{th})\), and \(V_{dep} \sim -0.8VDD\). While adjustments in the details of this choice are often made in the interest of optimization of processes for a particular product, we will assume here this approximate equality of turn-on voltages of the two transistor types for the sake of simplicity. Substituting this choice of \(V_{dep}\) into eq.6a, we find that for \(V_{th}\) small compared to VDD:

\[
V_{inv} \sim VDD/[Z_{pu}/Z_{pd}]^{1/2}
\]  
[eq.7]

In general it is desirable that the margins around the inverter threshold be approximately equal, i.e., that the inverter threshold, \(V_{inv}\), lie approximately midway between VDD and ground. We see from eq.7 that this criterion is met by a ratio of pullup Z to pulldown Z of approximately 4:1. We will see later that the choice of \(V_{dep} \sim VDD - V_{th}\) producing a ratio of 4:1 here, will lead to a balancing of performances in certain other important circuits.
Inverter Delay

A minimum requirement for an inverter is that it drive another identical to itself. Let us analyze the delay through a string of inverters of identical dimensions. This is the simplest case in which we can estimate performance. Inverters connected in this way are shown in Fig. 4a. We define the inverter ratio $k$ as the ratio of $Z$ of the pullups to $Z$ of the pulldowns. We will sometimes use the alternative "resistor with gate" pullup symbol, as in fig.4a, to clarify its functional purpose.

Let us assume that prior to $t = 0$, the voltage at the input of the first inverter is zero, and hence, the voltage output of the second inverter will be low. At time $t=0$, let us place a voltage equal to VDD on the input of the first inverter and follow the sequence of events which follows. The output of the first inverter, which leads to the gate of the second inverter, will initially be at VDD. Within approximately one transit time, the pulldown transistor of the first inverter will remove from this node an amount of charge equal to VDD times the gate capacitance of the pulldown of the second inverter. The pulldown transistor of the second inverter is now faced with the task of supplying a similar charge to the gate of the third inverter, to raise it to VDD. Since it can supply at most only $1/k$th of the current that can be supplied by the pulldown transistor, the delay in the second inverter stage is approximately $k$ times that of the first.

It is thus convenient to speak of the inverter pair delay which includes the delay for one lowgoing transition and one highgoing transition. This inverter pair delay is approximately $(1+k)$ times the transit time, as shown in figure 4a. The fact that the rising transition is slower than the falling transition by approximately the inverter transistors' geometry ratios is an inherent characteristic of any ratio type logic. It is not true of all logic families. For example, in families such as complementary MOS where there are both pMOS and nMOS devices on the same silicon chip and both types operate strictly as pulldown enhancement mode devices, any delay asymmetry is a function of the difference in mobilities of the p and n type charge carriers rather than of the transistor geometrical ratios.

Fig. 4b shows an inverter driving the inputs of several other inverters. In this case, for a fanout factor $f$, it is clear that in either the pullup or pulldown direction, the active device must supply $f$ times as much charge as it did in the case of driving a single input. In this case, the delay both in the up and downgoing directions is increased by approximately the factor $f$. In the case of the
downgoing transition, the delay is approximately $f$ times the transit time of the pulldown 
transistor, and in the case of the upgoing transition, the delay is approximately the inverter ratio $k$ 
times the fanout factor times the pulldown transit time.

In the discussions of transit time given earlier, it was assumed that both the depletion mode 
pullup device and the enhancement mode pulldown device were operating in the resistive region. 
It was also assumed that all capacitances were constant, and not a function of voltage. These 
conditions are not strictly met in the technology we are discussing. Delay calculations given in 
this text are based on a "switching model" where individual stages spend a small fraction of their 
time in the mid-range of voltages around $V_{\text{inv}}$. This assumption introduces a small error of the 
order of $1/G$. Because of these and other second order effects, the switching times actually 
observed vary somewhat from those derived.

Parasitic Effects

In integrated systems, capacitances of circuit nodes are due not only to the capacitance of gates 
connected to the nodes, but also include capacitances to ground of signal paths connected to the 
nodes and other stray capacitances. These other capacitances, sometimes called parasitic or stray 
capacitances, are not negligible. While gate capacitances are typically an order of magnitude 
greater per unit area than the capacitances of the signal paths, the signal paths are often much 
larger in area than the associated gate regions. Therefore, a substantial fraction of the delay 
encountered may be accounted for by stray capacitance rather than by the inherent properties of 
the active transistors. In the simplest case where the capacitance of a node is increased by the 
presence of parasitic area attached to the node, the delays can be accounted for by simply 
increasing the transit time by the ratio of the total capacitance to that of the gate of the transistor 
being driven. Time is required to supply charge not only to the gate itself but also to the 
parasitic capacitance.

There is one type of parasitic, however, which is not accounted for so simply. All MOS 
transistors have a parasitic capacitance between the drain edge of the gate and the drain node. 
This effect is shown schematically in figure 4c. In an inverter string, this capacitance will be 
charged in one direction for one polarity of input, and in the opposite direction for the opposite 
polarity input. Thus, on a gross scale its effect on the system is twice that of an equivalent 
parasitic capacitance to ground. Therefore, gate to drain capacitances should be approximately
Fig. 4a. Inverter Delay

For fanout of f:

down delay $\sim f \cdot \tau$

up delay $\sim k f \cdot \tau$

Fig. 4b. Fanout

Start:

$V2 = 0, V3 = VDD$

$Qg = 0, Qm = -CmVdd$

$Qwire = 0$

Finish:

$V2 = VDD, V3 = 0$

$Qg = CgVDD, Qm = +CmVDD$

$Qwire = CstrayVDD$

Fig. 4c. The Miller Effect

$\text{total effective input capacitance} = \frac{dQ}{dV} = \frac{(Q_{\text{finish}} - Q_{\text{start}})}{(V_{\text{finish}} - V_{\text{start}})} = C_g + 2C_m + Cstray$

(most press)
doubled, and added to the gate capacitance $C_g$ and the stray capacitances, to account for the total capacitance of the node and thus for the effective delay time of the inverter. The effective inverter pair delay then is equal to $\tau(1+k)C_{\text{total}}/C_g$.

**Driving Large Capacitive Loads**

As we have seen, the delay per inverter stage is multiplied by a fanout factor. The overall performance of a system may be seriously degraded if it contains any large fanouts, where one circuit within the system is required to drive a large capacitive load. As we shall see, this situation often occurs in the case of control drivers required to drive a large number of inputs to memory cells or logic function blocks. A similar and more serious problem is driving wires which go off the silicon chip to other chips or input/output devices. In such cases the ratio of the capacitance which must be driven to the inherent capacitance of a gate circuit on the chip is often many orders of magnitude, causing a serious delay and a degradation of system performance.

Consider how we may drive a capacitive load $C_L$ in the minimum possible time given that we are starting with a signal on the gate of an MOS transistor of capacitance $C_g$. Define the ratio of the load capacitance to the gate capacitance, $C_L/C_g$, as $Y$. It seems intuitively clear that the optimum way to drive a large capacitance is to use our elementary inverter to drive a larger inverter and that larger inverter to drive a still larger inverter until at some point the larger inverter is able to drive the load capacitance directly. Using an argument similar to the fanout argument it is clear that for one inverter to drive another inverter, where the second is larger in size by a factor of $f$, results in a delay $f\tau$ times the inherent inverter delay, $\tau$. If $N$ such stages are used, each larger than the previous by a factor $f$, then the total delay of the inverter chain is $Nf\tau$, where $f^N$ equals $Y$. Note that if we use a large factor $f$, we can get by with few stages, but each stage will have a long delay. If we use a smaller factor $f$, we can shorten the delay of each stage, but are required to use more stages. What value of $N$ minimizes the overall delay for a given $Y$?

We compute this value as follows:

Since $f^N = Y, \quad \ln(Y) = N \ln(f)$

Delay of one stage $= f\tau$.

Thus total delay is $\quad = Nf\tau = \ln(Y)[f/\ln(f)]\tau$  \[eq.8\]
Notice that the delay is always proportional to \( \ln(Y) \), a result of the exponential growth in successive stages of the driver. The multiplicative factor, \( f/\ln(f) \), is plotted as a function of \( f \) in figure 5, normalized to its minimum value (e). Total delay is minimized when each stage is larger than the previous one by a factor of e, the base of natural logarithms. Minimum total delay is the elementary inverter delay \( \tau \) times e times the natural logarithm of the ratio of the load capacitance to the elementary inverter capacitance.

\[
\text{Min. total delay} \sim re[\ln(C_L/C_g)] \quad [\text{eq.9}]
\]

Minimum delay through the driver is seldom the only design criterion. The relative time penalty introduced by the choice of other values of \( f \) can be read directly from figure 5.

**Space vs Time**

From the results of the sections on inverter delay, parasitic effects, and driving large capacitances, we see that areas and distances on the silicon surface trade off against delay times. For an inverter to drive another inverter some distance away, it must charge not only the gate capacitance of the succeeding inverter but also the capacitance to ground of the signal path connecting the two. Increasing the distance between the two inverters will therefore increase the inverter pair delay. This effect can be counterbalanced by increasing the area of the first inverter, so as to reduce the ratio of the load capacitance to the gate capacitance of the first inverter. But the delay of some previous driving stage is then increased. There is no way to get around the fact that transporting a signal from one node to another some distance away requires charging or discharging capacitance, and therefore takes time. Note that this is not a velocity of light limitation as is often the case outside the chip. The times are typically several orders of magnitude longer than those required for light to traverse the distances involved. To minimize both the time and space required to implement system functions, we will tend to use the smallest possible circuits and locate them in ways which tend to minimize the interconnection distances.

The results of a previous section can be used here to illustrate another interesting space vs time effect. Suppose that the minimum size transistors of an integrated system have a transit time \( \tau \) and gate capacitance \( C_g \). A minimum size transistor within the system produces a signal which is then passed through successively larger inverting logic stages and eventually drives a large
Fig. 5. Relative Time Penalty, $f^* \frac{f}{e^{\ln(f)}}$ vs Size Factor $f$.
capacitance $C_L$ with minimum total delay equal to $t_{min}$. With the passage of time, fabrication technology improves. We replace the system with one in which all circuit dimensions, including those vertical to the surface, are scaled down in size by dividing by a factor $\alpha$, and the values of $V_{dd}$ and $V_{th}$ are also scaled down by dividing by $\alpha$. The motivation for this scaling is clear: the new system may contain $\alpha^2$ as many circuits. As described in a later section, we will find that the transit times of the smallest circuits will now be $\tau' = \tau/\alpha$, and their gate capacitance will be $C_g' = C_g/\alpha$. The new ratio of load to minimum gate capacitance is: $Y' = \alpha Y$. Referring to equation 9., we find that the new minimum total delay, $t_{min}'$, to drive $C_L$ scales as follows:

$$t_{min}' = t_{min} \left(1/\alpha\right) \left[1 + (\ln\alpha/\ln Y)\right]$$

Therefore, as the inverters scale down and $\tau$ gets smaller, more inverting logic stages are required to obtain the minimum offchip delay. Thus the relative delay to the outside world becomes larger. However, the absolute delay becomes smaller.

Basic NAND and NOR Logic Circuits

NAND and NOR logic circuits may be constructed in integrated systems as simple expansions of the basic inverter circuit. The analysis of the behavior of these circuits, including their logic threshold voltages, transistor geometry ratios and time delays, is also a direct extension of the analysis of the basic inverter.

The circuit layout diagram of a two input NAND gate is shown in figure 6a. The layout is that of a basic inverter with an additional enhancement mode transistor in series with the pulldown transistor. NAND gates with more inputs may be constructed by simply adding more transistors in series with the pulldown path. The electrical circuit diagram, truth table and logic symbol for the two input NAND gate are shown in figure 6b. If either of the inputs A or B is a logic-0, the pulldown path is open and the output will be high, and therefore a logic-1. For the output to be driven low, to logic-0, both inputs must be high, at logic-1. The logic threshold voltage of this NAND gate is calculated in a similar manner to that of the basic inverter, except equation 7 is rewritten with the length of the pulldowns replaced with the sum of the lengths of the two pulldowns (assuming their widths are equal) as follows:
\[ V_{thNAND} \sim VDD/[(L_{pu}/W_{pu})/((L_{pd_a} + L_{pd_b})/W_{pd})]^{1/2} \]

This equation indicates that as pulldowns are added in series to form NAND gate inputs, the pullup length must be enlarged to hold the logic threshold voltage constant.

The logic threshold voltage of an n-input NAND gate, assuming all the pulldowns have equal geometries, is:

\[ V_{thNAND} \sim VDD/[(L_{pu}/W_{pu})/(nL_{pd}/W_{pd})]^{1/2} \]

As inputs are added and pullup length is increased, the delay time of the NAND gate is also correspondingly increased, both for rising and falling transitions.

\[ \tau_{NAND} \sim n\tau_{inv} \]

The circuit layout diagram of a two input NOR gate is shown in figure 6c. The layout is that of a basic inverter with an additional enhancement mode transistor in parallel with the pulldown transistor. Additional inputs may be constructed by simply placing more transistors in parallel with the pulldown path. The circuit diagram, truth table and logic symbol for the two input NOR gate are shown in figure 6d. If either of the inputs A or B is a logic-1, the pulldown path to ground is closed and the output will be low, and therefore a logic-0. For the output to be driven high, to logic-1, both inputs must be low, at logic-0. If one of its inputs is kept at logic-0, and the other swings between logic-0 and logic-1, the logic threshold voltage of the NOR gate is the same as that of a basic inverter of equal pullup to pulldown ratio. If this ratio were 4:1 to provide equal margins, then \( V_{thNOR} \sim VDD/2 \) with only one input active. However, if both pulldowns had equal geometries, and if both inputs were to move together between logic-0 and logic-1, \( V_{thNOR} \) would be reduced to \( \sim VDD/(8)^{1/2} \). The logic threshold voltage of an n-input NOR circuit decreases as a function of the number of active inputs (inputs moving together from logic-0 to logic-1). The delay time of the NOR gate with one input active is the same as that of an inverter of equal transistor geometries, except for added stray capacitance. Its delay time for falling transitions is decreased as more of its inputs are active.
Fig. 6a. NAND Gate  
[Top view of layout]

Fig. 6b. NAND Gate Circuit Diagram,  
Logic Symbol, Logic Function

Fig. 6c. NOR Gate  
[Top view of layout]

Fig. 6d. NOR Gate Circuit Diagram,  
Logic Symbol, Logic Function
Fig. 7a. Inverting Super Buffer

Fig. 7b. Non-Inverting Super Buffer
Super Buffers

As we have noted, ratio type logic suffers from an asymmetry in its ability to drive capacitive loads. This asymmetry results from the fact that the pullup transistor has of necessity less driving capability than the pulldown transistor. There are, however, methods for avoiding this asymmetry. Shown in figures 7a and 7b are circuits for inverting and non-inverting drivers which are approximately symmetrical in their capability of sourcing or sinking charge into a capacitive load. Drivers of this type are called super buffers.

Both types of super buffer are built using a depletion mode pullup transistor and an enhancement mode pulldown transistor, with a ratio of Z's of approximately 4:1 as in the basic inverter. However, the gate of the pullup transistor, rather than being tied to its source, is tied to a signal which is the complement of that driving the pulldown transistor.

When the pulldown transistor gate is at a high voltage, the pullup transistor gate will be approximately at ground, and the current through the super buffer will be similar to that through a standard inverter of the same size. However, when the gate of the pulldown transistor is put to zero, the gate of the pullup transistor will go rapidly to VDD since it is the only load on the output of the previous inverter, and the depletion mode transistor will be turned on at approximately twice the drive which it would experience if its gate were tied to its source. Since the current from a device in saturation goes approximately as the square of the gate voltage, the current sourcing capability of a super buffer is approximately four times that of a standard inverter. Hence, the current sourcing capability of its pullups is approximately equal to the current sinking capability of its pulldowns, and wave forms from super buffers driving capacitive loads are nearly symmetrical.

The effective delay time, \( \tau \), of super buffers is thus reduced to approximately the same value for highgoing and lowgoing wave forms. Needless to say, when large capacitive loads are to be driven, super buffers are universally used. The arguments used in the last section to determine how many stages are used to drive a large capacitive load from a small source apply directly to super buffers. For that reason we have not explicitly indicated an inverter ratio \( k \) in that section.
A Closer Look at the Electrical Parameters

Up to this point we have talked in very simple terms about the properties of the MOS transistors. They have a capacitance associated with their gate input and a transit time for electrons to move from the source to the drain. We have given simple expressions for the drain to source current. For very low $V_{ds}$, the MOS transistor's drain to source path acts as a resistor whose conductance is directly proportional to the gate voltage above threshold, as given in equation 3. For values of $V_{ds}$ larger than $V_{gs} - V_{th}$, the device acts as a current source, with a current proportional to $(V_{gs} - V_{th})^2$, as given in equation 5. As $V_{ds}$ passes through the intermediate range between these two extremes, there is a smooth transition between the two types of behavior\(^1\), as given in the following equation:

$$I_{ds} = Q/\tau = \mu C_g (V_{gs} - V_{th}) V_{ds} - (V_{ds})^2/2/L^2$$  \[eq.10\]

Figure 9a plots $I_{ds}$ vs $V_{ds}$ summarizing the various regions of MOS transistor operation.

There is another electrical characteristic we may occasionally have to take into account. The threshold voltage of an MOS transistor is not a constant, but varies slightly as a function of the voltage between the source terminal of the transistor and the silicon substrate, which is usually at ground. This so called body effect is illustrated in figure 9b. If the source to bulk (substrate) voltage, $V_{sb}$, equals zero, then $V_{th}$ is at its minimum value of approximately 0.2 VDD. As $V_{sb}$ is increased, $V_{th}$ increases slightly.

For enhancement mode transistors fabricated using typical processes, $V_{th}$ reaches a maximum value of about 0.3 VDD when $V_{sb}$ is increased to ~ VDD. The value of the depletion mode transistor threshold, $V_{dep}$, is similarly affected, ranging from about -0.8 VDD to -0.7 VDD as $V_{sb}$ is raised from zero to VDD volts. As shown in figure 9b, it is possible to insert a fixed bias voltage between the circuit ground and the substrate, rather than just connect them. Such a substrate bias provides an electrical mechanism for setting the threshold to an appropriate value.
Fig. 9a. Summary of MOS Transistor Characteristics

(a) Cutoff Region:
Vgs < Vth, Ids = 0

(b) Saturation Region:
Vgs ≥ Vth, Vds sufficiently high so
Vgd < Vth. (i.e.: Vds ≥ Vgs - Vth)
MOSFET acts as current source, with
Ids proportional to \( (Vgs - Vth)^2 \)

(c) Resistive Region:
Vgs ≥ Vth, Vds sufficiently low so
Vgd ≥ Vth. (i.e.: Vds < Vgs - Vth)
MOSFET acts as resistor, with resistance
inversely proportional to \( (Vgs - Vth) \)

Fig. 9b. The Body Effect
Depletion Mode vs Enhancement Mode Pullups

With its gate tied to VDD, an enhancement mode transistor will be on for all $V_{ds} > V_{th}$, and thus can be used for a pullup device in inverting logic circuits. Early MOS processes used pullup devices of exactly this type.

In this section we will make a comparison of the rising transients of the two types of pullup circuits. As noted earlier, rising transients in ratio type logic are usually slower than falling transients, and thus rising transients generally have greater impact on system performance. In the simplest cases, this asymmetry in the transients results from the current sourcing capability of the pullup transistor being less than that of its pulldown counterpart. The simple intuitive arguments given earlier are quite adequate for making estimates of system performance in most cases. However, there are situations in which the transient time may be much longer than a naive estimate would indicate. The rising transient of the enhancement mode pullup is one of these.

A depletion mode pullup transistor feeding a capacitive load is shown schematically in figure 10a. Since $V_{gs} \geq V_{th}$ and $V_{gd} \geq V_{th}$, the pullup transistor is in the resistive region. The final stages of the rising transient are given by the following exponential:

$$V(t) = VDD\left[1 - e^{-t/(RC_L)}\right]$$

For an inverter ratio $k$, pulldown transit time $\tau$, and gate capacitance $C_g$, the time-constant of the rising transient is given by:

$$RC_L = k\tau C_L / C_g$$

A somewhat more complicated situation is presented by an enhancement mode transistor sourcing charge into a capacitive load. This situation is shown schematically in Fig. 10b. Note that since $V_{gd} = 0$, the transistor is in saturation whenever $V_{gs} > V_{th}$. The problem with sourcing charge from the enhancement mode transistor is that as the voltage at the output node gets closer and closer to one threshold below VDD, the amount of current provided by the enhancement mode transistor decreases rapidly.
The dependence of the enhancement mode pullup current, $I_{ds}$, upon output voltage, $V$, is given in equation 11:

$$Q = -\frac{\varepsilon W L}{D} [(VDD - V_{th}) - V]$$

$$\tau = \frac{2L^2/\mu}{(VDD - V_{th}) - V}$$

$$I_{ds} = -\frac{Q}{\tau} = \frac{\mu W}{2LD} [(VDD - V_{th}) - V]^2$$

[eq.11]

The fact that the pullup current decreases as the output voltage nears its maximum value causes the rising transient from such a circuit to be of qualitatively different form than that of a depletion mode pullup. Equating $I_{ds} = C_L dV/dt$ with the expression in equation 11, and then solving for $V(t)$, we find the rising voltage transient, for large $t$:

$$V(t) = VDD - V_{th}' - C_L \frac{L D}{\mu \varepsilon W t}$$

[eq.12]

Note that in this configuration, the threshold voltage $V_{th}'$ of the pullup is near its maximum value as $V(t)$ rises towards $VDD$, due to the body effect.

A comparison of the rising transients of the preceding two circuits, assuming the same load capacitance and the same pullup source current at zero output voltage, is shown in Fig. 10c. The rising transient for the depletion mode pullup transistor is crisp and converges rapidly towards $VDD$. However, the rising transient for the enhancement mode pullup transistor, while starting rapidly, lags far behind and within the expected time response of the system, never even comes close to one threshold below $VDD$. Even for very large $t$, $V(t) < VDD - V_{th}'$.

The practical effect of this property of enhancement mode transistors is that circuits designed to work from the output of such a circuit should be designed with an inverter threshold $V_{inv}$ considerably lower than that of circuits designed to work with the output of a depletion mode pullup circuit. In order to obtain equal inverter margins without sacrificing performance, we will normally use depletion mode pullups.
Fig. 10a. Depletion Mode
MOSFET Pulling Up
Capacitive Load

Fig. 10b. Enhancement Mode
MOSFET Pulling Up
Capacitive Load

Fig. 10c. Comparisons of Rising Transients for the Two Types of Pullups
Delays in Another Form of Logic Circuitry

Enhancement mode transistors, when used in small numbers and driving small capacitive loads, may often be used as switches in circuits of simple topology to provide logic signal steering functions of much greater complexity than could be easily achieved in ratio type inverting logic. These circuits are reminiscent of relay switching logic, and transistors used in this way are referred to as "pass transistors" or "transmission gates". Example circuits using this type of design are given in Chapter 3. A particularly interesting example is the Manchester carry chain\(^{4a,b}\), used for propagating carry signals in parallel adders. In each stage of the adder a carry propagate signal is derived from the two input variables to the adder, and if it is desired to propagate the carry, this propagate signal is applied to the gate of an enhancement mode pass transistor. The source of the transistor is carry-in to the present stage, and the drain of the transistor is carry-out to the next stage. In this way, a carry can be propagated from less to more significant stages of the adder without inserting a full inverter delay between stages. The circuit is shown schematically in Fig. 11a.

The delay through such a circuit does not involve inverter delays but is of an entirely different sort. A voltage along the chain divides into \( V_{ds} \) across each pass transistor. Thus \( V_{ds} \) is usually low, and the pass transistors operate primarily in the resistive region. We can think of each transistor as a series resistance in the carry path, and a capacitance to ground formed by the gate to channel capacitance of each transistor, and the strays associated with the source, drain, and connections with the following stage. An abstraction of the electrical representation is shown in Fig. 11b. The minimum value of \( R \) is the turned on resistance of each enhancement mode pass transistor, while the minimum value of \( C \) is the capacitance from gate to channel of the pass transistor. Strays will increase both values, especially that of \( C \). The response at the node labelled \( V_2 \) with respect to time is given in eq. 13. In the limit as the number of sections in the network becomes large, eq. 13 reduces to the differential form shown in eq. 14 where \( R \) and \( C \) are now the resistance and capacitance per unit length, respectively.

\[
C \frac{dV_2}{dt} = \frac{(V_1 - V_2) - (V_2 - V_3)}{R} \quad [eq. 13]
\]

\[
\frac{RC}{dV}{dt} = \frac{d^2V}{dx^2} \quad [eq. 14]
\]
Equation 14 is the well-known diffusion equation, and while its solutions are complex, in general the time required for a transient to propagate a distance $x$ in such a system is proportional to $x^2$. One can see qualitatively that this might be so. Doubling the number of sections in such a network doubles both the resistance and the capacitance, and therefore causes the time required for the system to respond to increase by a factor of approximately four. The response of a system of $n$ stages to a step function input is shown in Fig. 11c.

If we add one more pass transistor to such a chain of $n$ pass transistors, the added delay through the chain is small for small $n$, but very large for large $n$. Therefore, it is highly desirable to group the pass transistors used for steering, multiplexing, and carry-chain type logic into short sections and interpose inverting logic between these sections. This approach applied to the carry chain is shown in figure 11d. The delay through a section of $n$ pass transistors is proportional to $RCn^2$. Thus the total delay is $\sim RCn^2$ plus the delay through the inverter $\tau_{\text{inv}}$. The average delay per stage is given in eq. 15. To minimize the delay per stage, chose $n$ such that the delay through $n$ pass transistors equals the inverter delay:

\[
\text{Total delay} \sim RCn^2 + \tau_{\text{inv}},
\]

\[
\text{Average delay/stage} \sim RC + \tau_{\text{inv}}/n
\]  \[\text{[eq. 15]}\]

\[
\text{Min. delay when: } RCn^2 \sim \tau_{\text{inv}}
\]

Since logic done by steering signals with pass transistors does not require static power dissipation, a generalization of this result may be formulated. It pays to put as much logic into steering type circuits as possible until there are enough pass transistors to delay the signal by approximately one inverting logic delay. At this point, the level of the signal can be restored by an inverting logic stage.

The pass transistor has another important advantage over an inverting logic stage. When used to control or steer a logic signal, the pass transistor has only an input, control, and output connections. A NAND or NOR logic gate implementing the same function, in addition to containing two more transistors and thus occupying more area, also requires VDD and GND connections. As a result, the topology of interconnection of pass transistor circuits is far simpler than that of inverting logic circuits. This topological simplicity of pass transistor control gates is an important factor in the system design concepts developed in later chapters.
Fig. 11a. Pass Transistor Chain Propagating a Carry Signal

Fig. 11b. Equivalent Circuit

Fig. 11c. Response to Step Function Input

Fig. 11d. Minimizing Delay by Interposing Inverters
Fig. 12a. Inverters Coupled by Pass Transistor

\[ \frac{Z_{pu1}}{Z_{pd1}} = 4 \]

Fig. 12b. For \( V_{out2} = V_{out1} \), \( \frac{Z_{pu2}}{Z_{pd2}} = 8 \)
Pullup/Pulldown Ratios for Inverting Logic Coupled by Pass Transistors

Earlier we found that when an inverting logic stage directly drives another such stage, a pullup to pulldown ratio \( Z_{pu}/Z_{pd} = (L_{pu}/W_{pu})/(L_{pd}/W_{pd}) \) of 4:1 yields equal inverter margins, and also provides an output sufficiently less than \( V_{th} \) for an input equal to VDD. Rather than coupling inverting logic stages directly, we often couple them with pass transistors for the reasons developed in the preceding section, thus affecting the required pullup to pulldown ratio.

Figure 12a shows two inverters connected through a pass transistor. If the output of the first inverter nears VDD, the input of the second inverter can rise at most to \( (VDD - V_{thp}) \), where \( V_{thp} \) is the threshold of the pass transistor. Why does this effect occur? Consider the following: The output of the first inverter is at or above \( (VDD - V_{thp}) \), the pass transistor gate is at zero volts, and the input gate of the second inverter is also at zero volts. The pass transistor's gate voltage is now driven quickly to VDD, turning on the pass transistor. As current flows through the pass transistor, from drain to source, the input gate voltage of the second inverter rises and the gate to source voltage of the pass transistor falls. When the gate voltage of the second inverter has risen to \( (VDD - V_{thp}) \), the pass transistor's gate to source voltage has fallen to its threshold value, and the pass transistor will switch off.

If the second inverter is to have its output driven as low with an input of \( (VDD - V_{thp}) \) as would a standard inverter with an input of VDD, then the second inverter must have a pullup to pulldown ratio larger than 4:1. This larger ratio is calculated as follows: With inputs near VDD, the pullups of inverters are in saturation, and the pulldowns are in the resistive region. Figure 12b shows equivalent circuits for two inverters. VDD is input to one, and \( (VDD - V_{thp}) \) to the other. For the output voltages of the two inverters to be equal under these conditions, \( I_1R_1 \) must equal \( I_2R_2 \). Referring to equations 3a and 5, we find:

\[
(Z_{pu1}/Z_{pd1})(VDD - V_{th}) = (Z_{pu2}/Z_{pd2})(VDD - V_{th} - V_{thp})
\]

Since \( V_{th} \) of the pulldowns is approximately 0.2VDD, and \( V_{thp} \) of the pass transistor is approximately 0.3VDD due to the body effect, then \( Z_{pu2}/Z_{pd2} \sim Z_{pu1}/Z_{pd1} \). Thus a ratio of \( (L_{pu}/W_{pu})/(L_{pd}/W_{pd}) = 8 \) is usually used for inverting logic stages placed as level restorers between sections of pass transistor logic.

[ Ch1.: Sect.1 ] < Conway > newmos2.visi: July 1, 1978 3:13 PM
Transit Times and Clock Periods

In chapter 3 we will develop a system design methodology in which we will be able to construct and estimate the performance of arbitrarily complex digital systems, using only the basic circuit forms presented in the preceding sections. The basic system building block in the design methodology is a register to register transfer through combinational logic, implemented with pass transistors and inverting logic stages. Using the basic ideas already presented, we may anticipate the results of that chapter in order to estimate the maximum clocking frequency of such systems.

The design methodology uses a two-phase non-overlapping clock scheme. During the first clock phase, data passes from one register, through combinational logic stages and pass transistors to a second register. During the second clock phase, data passes from the second register through still more logic and pass transistors to a third (or possibly back to the first) register. The data storage registers are implemented by using charge stored on the input gates of inverting logic stages, the charge being isolated by pass transistors controlled by clock signals, as described in chapter 3.

Since pass transistors are used to connect inverting logic stages, inverter ratios of $k \sim 8$ are required. If the combinational logic between registers is implemented using only pass transistors, and if the delays through the pass transistors have been carefully matched to those of the inverting logic stages, the total delay will be twice that of the simple $k = 8$ inverter. In the absence of strays, the $k = 8$ inverters have a maximum delay (in the case of the output rising towards VDD) of $8\tau$, and hence a minimum of $16\tau$ must be allowed for the inverter plus logic delay. However, in most designs the stray capacitance is at least equal to that inherent in the circuit. Thus the minimum time required for one such operation is $\sim 30\tau$. Control lines to the combinational logic and pass transistors each typically drive the gates of 10 to 30 transistors. Even when using a super buffer driver, the delay introduced by this fan out is at least the minimum driving time for a capacitive load. With $Y = 30$, this time is $\sim 9\tau$. To this we must add an $8\tau$ inverter delay for operation of the drivers.

Thus the total time for one clock phase is $\sim 50\tau$. Since two clock phases are required per cycle, a minimum clocking period of $\sim 100\tau$ is required for system designed in this way. In 1978, $\tau \sim 0.3$ nanoseconds, and clocking periods of 30 to 50 ns are achievable in carefully structured integrated systems where successive stages are in close physical proximity. If it is necessary to communicate data over long distances, longer periods are required.
Properties of Cross-Coupled Circuits

In many control sequencing and data storage applications, memory cells and registers are built using two inverters driving each other, as shown in figure 13a. This circuit can be set in either the state where \( V_1 \) is high and \( V_2 \) is low, or in the state where \( V_1 \) is low and \( V_2 \) is high. In either case, the condition is stable and will not change to the other condition unless it is forced there through some external means. The detailed methods of setting such cross-coupled circuits into one state or another will be discussed in detail later. However, it is important at the present time to understand the time evolution of signals impressed upon cross-coupled circuits, since they exhibit properties different from circuits not having a feedback path from their output to an input.

We have seen that there exists a voltage at which the output of an inverter is approximately equal to its input voltage. If a cross-coupled circuit is inadvertently placed in a situation where its input voltage is equal to this value, then an unstable equilibrium condition is created where voltages \( V_1 \) and \( V_2 \) are equal. Since the net current flowing onto either gate is now zero, there is no forcing function driving the system to any voltage other than this equilibrium one, and the circuit can stay in this condition for an indefinite period. However, if either voltage changes, even very slightly, the circuit will leave this unstable equilibrium. For example, if the voltage \( V_1 \) is increased from its unstable equilibrium value \( V_{\text{inv}} \) by a slight amount, this will in time cause a lowering of voltage \( V_2 \), as net current flows from gate 1. This lowering of \( V_2 \) will at some later time cause \( V_1 \) to increase further. As time goes on, the circuit will feedback on itself until it rests in a stable equilibrium state.

The possibility of such unstable equilibria in cross-coupled circuits has important system implications\(^2\), as we will later see. For this reason, we will make a fairly detailed analysis of this circuit's behavior near the metastable state. While it is not essential that the reader follow all the details of the analysis, the final result should be studied carefully. The time constant of the final result depends in detail on the regions of operation of the transistors near the metastable state, as given in the following analysis. However, the exponential form of the result follows simply from the fact that the forcing function pushing the voltage away from the metastable point is proportional to the voltage's distance away from that point. This general behavior is characteristic of bistable storage elements in any technology. However, more complex waveforms are observed in logic families having more than one time constant per stage.

The time evolution of this process can be traced as follows. At the unstable equilibrium, the current in the pullups equals that in the pulldowns, and is some constant, \( k_1 \), times \((V_{\text{inv}} - V_{\text{th}})^2\). If \( V_1 \) is then changed by some small \( \Delta V_1 \) to \( V_{\text{init}} \), \( I_{\text{pu2}} \) remains constant but \( I_{\text{pd2}} \) changes immediately, producing a non-zero \( I_{g1} \):

\[
I_{g1} = I_{\text{pu2}} - I_{\text{pd2}} = k_1[(V_{\text{inv}} - V_{\text{th}})^2 - (V_{\text{inv}} + \Delta V_1 - V_{\text{th}})^2]
\]

For small \( \Delta V_1 \), \( I_{g1} = -2k_1(V_{\text{inv}} - V_{\text{th}})\Delta V_1 \). More precisely, since \( I_{g1} = \text{function}(V_1, V_2) \), then near \( V_{\text{inv}} \):

\[
\frac{\partial I_{g1}}{\partial V_1} = -2k_1(V_{\text{inv}} - V_{\text{th}})
\]

Noting that the pullups are not quite in saturation, but are in the resistive region, and:

\[
\frac{\partial I_{g1}}{\partial V_2} = -1/R_{\text{pu}},
\]

where \( R_{\text{pu}} \) = effective resistance of the pullup near \( V_{\text{inv}} \). Noting that \( I_{g1} = C_g \frac{dV_{g2}}{dt} \), we find that:

\[
\frac{dI_{g1}}{dt} = -2k_1(V_{\text{inv}} - V_{\text{th}}) \frac{dV_1}{dt} - (1/R_{\text{pu}}) \frac{dV_2}{dt} = C_g \frac{d^2V_2}{dt^2}
\]

Evaluating the constants in this equation yields \(-k_1(V_{\text{inv}} - V_{\text{th}}) = C_g/\tau_0 \) where \( \tau_0 \) is the saturation transit time of the pulldowns for \( t \) near zero. Assume a pullup/pulldown \( Z \) ratio of 4:1, and consider the operating conditions near \( t = 0 \). Evaluating the effective resistance of the pullups in terms of the parameters of the pulldowns yields \( 1/R_{\text{pu}} \sim C_g/\tau_0 \).

Therefore:

\[
-(2/\tau_0) \frac{dV_1}{dt} - (1/\tau_0) \frac{dV_2}{dt} = \frac{d^2V_2}{dt^2}
\]

Similarly:

\[
-(2/\tau_0) \frac{dV_2}{dt} - (1/\tau_0) \frac{dV_1}{dt} = \frac{d^2V_1}{dt^2}
\]
Fig. 13a. Cross Coupled Inverters

Fig. 13b. $V(t)$ for Cross Coupled Inverters

(mos13ab press)
Near time \( t = 0 \), \( \frac{dV_1}{dt} \) approximately equals \( -\frac{dV_2}{dt} \), and therefore:

\[
\frac{d^2V_1}{dt^2} = -(1/\tau_o) \frac{dV_2}{dt} = (1/\tau_o)^2 V_1 + \text{const.} \quad \text{[eq. 16a]}
\]

The solution to eq. 16a is an exponential diverging from the equilibrium voltage \( V_{\text{inv}} \), with a time constant \( \tau_o/2 \) equal to \( 1/2 \) the pulldown delay time. Note that the solution given in eq. 16b satisfies the conditions that \( V(0) = V_{\text{init}} \) and that \( V(t) \) is constant, if \( V_{\text{init}} = V_{\text{inv}} \):

\[
V_1(t) = V_{\text{inv}} + (V_{\text{init}} - V_{\text{inv}}) e^{t/\tau_o} \quad \text{[eq. 16b]}
\]

The above analysis applies to cross coupled circuits in the absence of noise. Noise unavoidably present in the circuit spreads the input voltage into a band from which such an unstable equilibrium can statistically be initiated. The width of this band is equal to the noise amplitude. Any timing condition which causes the input voltage to settle in this band has some probability of causing a balanced condition, from which the circuit may require an arbitrarily long time to recover. The time evolution of such a system is shown in Fig. 13b, for several initial voltages near \( V_{\text{inv}} \). The time for the cross-coupled system to reach one of its equilibria is thus logarithmic in the displacement from \( V_{\text{inv}} \), and is given approximately by eq. 16c:

\[
t \sim \tau_o \ln[V_{\text{inv}}/(V_{\text{init}} - V_{\text{inv}})] \quad \text{[eq. 16c]}
\]
A Fluid Model for Visualizing MOS Transistor Behavior

[Section Contributed by Carlo H. Sequin, U. C. Berkeley]

When designing circuits and systems, it is often useful to have some method for visualizing the physical behavior of the devices used as basic building blocks. This section develops such a method for the MOS transistor. Some readers of this text may be unfamiliar with semiconductor device physics, and would have difficulty visualizing what is going on inside an active semiconductor device, if device behavior were described in purely analytical terms. However, it is possible to construct a simple but very effective model of the behavior of certain charge controlled devices, such as MOS transistors, charge coupled devices (CCD's), and bucket brigade devices (BBD's)\(^8\), without referring to the details of device physics.

This model will be developed using two basic ideas: We think of electrical charge as though it were a fluid, and we mentally map the relevant electrical potentials into the geometry of a "container" in which the charge is free to move around. One can then apply one's intuitive understanding of, say, water in buckets of various shapes towards a visualization of what is going on inside the devices. Often a design guided by a good intuitive understanding of how a fluid would behave in the designed structure may show superior performance over designs based on complicated but possibly inadequate two-dimensional analytical modelling.

The MOS Capacitor

The basic element of MOS transistors or charge transfer devices is the MOS capacitor. The notions of a fluid model will first be introduced using this elementary building block.

In physical space an MOS capacitor is a sandwich structure of a metal or polysilicon electrode on a thin insulator on the surface of a silicon crystal (fig. fm-1a). A suitable voltage applied to the electrode, i.e. positive for a p-type silicon substrate as used in nMOS, will repel the majority carriers in the substrate under the electrode, generating a depletion region which is at first free of any mobile charge carriers. Minority carriers, in this case electrons, can be injected electrically into this area, or generated by incident light, and subsequently stored underneath the MOS electrode. Applying the notions of a fluid model, the same situation can be described as follows:

The positive voltage applied to the MOS electrode generates a pocket in the surface potential of the silicon substrate. This can be visualized as a container, where the shape of the container is
defined by the potential along the silicon surface, as plotted by the dashed line in figure fm-1b. Note that in fig. fm-1b, increasing positive potential is plotted in the downward direction. The presence of minority charge carriers in an inversion layer changes the surface potential: an increase in this charge decreases the positive surface potential under the MOS electrode. The potential profile in the presence of inversion charge is indicated by the solid line in fig. fm-1b. The area between the dashed and solid lines in fig. fm-1b is hatched to indicate the presence of this charge. This representation shows charge sitting at the bottom of the container, just as a fluid would reside in a bucket. Of course the surface of the fluid (solid line) must be level in an equilibrium condition; if it were not, electrons would move under the influence of the potential difference until a constant surface potential has been established.

This model allows one to visualize easily the amount of charge present (hatched area), the fact that the charge tends to sit in the deepest part of the potential well, and the fact that the capacity of the bucket is finite and dependent upon the applied electrode voltage. The higher this voltage, the deeper the bottom of the bucket and the more charge that can be stored. It should be kept in mind that this fluid model differs from the physical reality in so far as in reality the minority carriers in the inversion layer reside directly at the silicon surface.

The MOS Transistor

The same kind of model can be used to describe MOS transistor behavior. Figure fm-2a shows the physical cross section through an MOS transistor. Source and drain diffusions have been added to the simple MOS capacitor. For the moment we consider these two diffusions to be connected to two identical voltage sources, $V_{sb} = V_{db}$, which thus define the potential of the source and drain regions.

In the potential plot these diffusions are represented by exceedingly deep buckets, filled with charge carriers up to the levels of the source and drain region potentials. Whether the MOS transistor is conducting, or is isolating the two diffused regions from one another, now depends on the potential underneath the MOS gate electrode. If the applied gate potential is chosen so that the potential underneath is less than $V_{sb}$, then there exists a potential barrier between source and drain regions (case 1 and 2 in fig. fm-2b). However, if the potential of an "empty bucket" under the gate electrode would be higher than $V_{sb}$, then the transistor is turned on (case 4 and 5). Of course, in cases 4 and 5, carriers from the source and drain regions will spill underneath the gate electrode so that a uniform surface potential exists throughout the whole transistor. The
conductivity of the channel area depends on the thickness of the inversion layer, which can readily be visualized in fig. fm-2b. Channel conductivity goes to zero at the turn-on threshold of the transistor (case 3), when the "empty bucket" potential under the gate electrode is equal to the source and drain potential. Thus, the region under the gate can be viewed as a movable barrier of variable height which controls the flow of charge between the source and drain areas.

The same model enables us to visualize what happens when source and drain regions are biased to different potentials, as is usually the case in normal operation of MOS transistors. Figure fm-3a again shows a physical cross section through an MOS transistor, as a reference for the following figures. Figure fm-3b reviews the case of equal source and drain potentials with the channel turned on fairly strongly, thus readily allowing charge to move between source and drain. Figure fm-3c shows the situation when a small voltage difference, ΔV, has been applied between source and drain. Since the potential difference is maintained by external voltage sources, electrons will be forced to move from source to drain under the influence of the potential gradient, just as a liquid would flow from the higher to the lower level.

As the potential difference between source and drain is made larger, the variation in the "depth" of the fluid along the channel becomes significant (fig. fm-3d). Continuity in the fluid requires that the charge move faster in the areas where the layer is thinner. This implies that the potential increases more rapidly closer to the drain region. With increasing drain potential the amount of charge flowing from source to drain per unit time increases, since the product of charge layer depth and local gradient increases. However, there is a limit. Once the drain potential exceeds the empty channel potential the rate of charge flow will be limited by the drain-side edge of the barrier under the gate electrode. The MOS transistor has now reached saturation (fig. fm-3e). The drain current density now is determined by the potential difference between the source and the empty channel and by the length of the channel (or the width of the barrier over which the charge has to flow), and is to first order independent of the drain voltage V_{db}.

Even in simple transistor circuits the above fluid model helps one quickly develop a feeling for device and circuit operation. However, the real power of this intuitive model emerges when it is applied to complex structures where closed form solutions describing charge motion can no longer be found. The empty potential under the various electrodes can first be plotted as in the above examples, and the flow of charge then visualized using the analogy to the behavior of a fluid.
Fig. fm3a:

MOS Transistor

p-type Si

Fig. fm3b:

\[ V_{sb} = V_{db} \]

Fig. fm3c:

\[ V_{db} = V_{sb} + \Delta V \]

Fig. fm3d:

\[ V_{db} > V_{sb} \]

Fig. fm3e:

\[ V_{db} > V_{gb} \cdot V_{th} \]
Effects of Scaling Down the Dimensions of MOS Circuits and Systems

This section examines the effects on major system parameters resulting from scaling down all dimensions of an integrated system, including those vertical to the surface, by dividing them by a constant factor $\alpha$. The voltage is likewise scaled down by dividing by the same constant factor $\alpha$. Using this convention, all electric fields in the circuit will remain constant. Thus many non-linear factors affecting performance will not change as they would if a more complex scaling were used.

Figure 14a. shows a MOSFET of dimensions $L$, $W$, $D$, with a $(V_{gs} - V_{th}) = V$. Figure 14b. shows a MOSFET similar to that in figure 14a., but of dimensions $L' = L/\alpha$, $W' = W/\alpha$, $D' = D/\alpha$, and $V' = V/\alpha$. Refer to equations 1., 2., and 3. From these equations we will find that as the scale down factor $\alpha$ is increased, the transit time, the gate capacitance, and drain to source current of every individual transistor in the system scale down proportionally, as follows:

$$\tau \propto L^2/V, \quad \tau' / \tau = [(L/\alpha)^2/(V/\alpha)]/[L^2/V], \quad \text{therefore,} \quad \tau' = \tau / \alpha$$

$$C \propto LW/D, \quad C'/C = (L/\alpha)(W/\alpha)/(D/\alpha)/[LW/D], \quad \text{and} \quad C' = C/\alpha$$

$$I \propto WV^2/LD, \quad I'/I = [(WV^2/\alpha^3)/(LD/\alpha^2)]/[WV^2/LD], \quad \text{and} \quad I' = I/\alpha$$

*Switching power, $P_{sw}$, is the energy stored on the capacitance of a given device divided by the clock period, or time between successive charging and discharging of the capacitance. A system’s clock period is proportional to the $\tau$ of its smallest devices. As devices are made smaller and faster, the clock period is proportionally shortened. Also, the dc power, $P_{dc}$, dissipated by any static circuit equals I times V. Therefore, $P_{sw}$ and $P_{dc}$ scale as follows:*

$$P_{sw} \propto CV^2/\tau \propto WV^3/DL, \quad \text{and} \quad P'_{sw} = P_{sw}/\alpha^2$$

$$P_{dc} = IV, \quad \text{and} \quad P'_{dc} = P_{dc}/\alpha^2$$
Both the switching power and static power per device scale down as $1/\alpha^2$. The average dc power for most systems can be approximated by adding the total $P_{sw}$ to one-half of the dc power which would result if all level restoring logic pulldowns were turned on. The contribution of pass transistor logic to the average dc power drawn by the system is due to the switching power consumed by the driving circuits which charge and discharge the pass transistor control gates.

The switching energy per device, $E_{sw}$, is an important metric of device performance. It is equal to the power consumed by the device at maximum clock frequency multiplied by the device delay, and scales down as follows:

$$E_{sw} \propto CV^2, \quad \text{and} \quad E_{sw} = E_{sw}/\alpha^3$$

The following table summarizes values of the important system parameters for current technology, and for a future technology near the limits imposed by physical law:

<table>
<thead>
<tr>
<th></th>
<th>1978</th>
<th>19XX</th>
</tr>
</thead>
<tbody>
<tr>
<td>Minimum Feature Size:</td>
<td>6 $\mu$m</td>
<td>0.3 $\mu$m</td>
</tr>
<tr>
<td>$\tau$</td>
<td>0.3 to 1 ns</td>
<td>~0.02 ns</td>
</tr>
<tr>
<td>$E_{sw}$</td>
<td>$\sim 10^{-12} J$</td>
<td>$\sim 2 \times 10^{-16} J$</td>
</tr>
<tr>
<td>System Clock Period:</td>
<td>~30 to 50 ns</td>
<td>~2 to 4 ns</td>
</tr>
</tbody>
</table>

[see earlier section]

A more detailed plot of the channel conductance of an MOS transistor near the threshold voltage is shown in figure 15. Below the nominal threshold, the conductance (1/R) is not in reality zero, but depends on gate voltage and temperature as follows:

$$1/R \propto e^{(V_{gs} - V_{th})/(kT/q)}$$

where $T$ is the absolute temperature, $q$ is the charge on the electron, and $k$ is Boltzmann’s constant. At room temperature, $kT/q \sim 0.025$ volts. At present threshold voltages, as in the
Fig. 14a. MOSFET, 1978

Fig. 14b. MOSFET Scaled Down by Alpha, 19XX

Fig. 15. Conductance as a Function of Threshold Voltage
rightmost curve in figure 15., an off device is below threshold by perhaps 20 kT/q, i.e. by about 0.5 volts, and its conductance is decreased by a factor of the order of ten million. Said another way, if the device is used as a pass transistor, a quantity of charge which takes a time $\tau$ to pass through the on device, will take a time on the order of $10^7\tau$ to "leak" through the off device.

The use of pass transistors switches to isolate and "dynamically store" charge on circuit nodes is common in many memory applications using 1978 transistor dimensions. However, if the threshold voltage is scaled down by a factor of perhaps 5, as shown in the leftmost curve in figure 15., then an off transistor is only 4kT/q below threshold. Therefore, its conductance when "off" is only a factor of 100 or so less than when it is "on". For such relatively large values of subthreshold conductance, charge stored dynamically on a circuit node by the transistor when "on" will safely remain on that node for only a few system clock periods. The charge will not remain on the node for a very large number of periods as it does in present memory devices using this technique. One way of possibly coping with this problem, as device dimensions and threshold voltages are scaled down, is to reduce the temperature of device operation.

Suppose we scale down an entire integrated system by a scale down factor of $\alpha = 10$. The resulting system will have one hundred times the number of circuits per unit area. The total power per unit area remains constant. All voltages in the system are reduced by the factor of 10. The current per unit area is increased by a factor of 10. The time delay per stage is decreased by a factor of 10. Therefore, the power-delay product decreases by a factor of 1000.

This is a rather attractive scaling in all ways except for the current density. The delivery of the required average dc current presents an important obstacle to scaling. This current must be carried to the various circuits in the system on metal conductors, in order that the voltage drop from the off-chip source to the on-chip subsystems will not be excessive. Metal paths have an upper current density limit imposed by a phenomenon called metal migration, discussed further in chapter 2. Many metal paths in today's integrated circuits are already operated near their current density limit. As the above type of scaling is applied to a system, the conductors get narrower, but still deliver the same current on the average to the circuits supplied by them.

Therefore, it will be necessary to find ways of decreasing system current requirements to approximately a constant current per unit area relative to the present current densities. In n-channel silicon gate technology, this objective can be partially achieved by using pass transistor...
logic in as many places as possible and avoiding restoring logic except where it is absolutely necessary. Numerous examples of this sort of design are given later in this text. This design approach also has the advantages of tending to minimize delay per unit function and to maximize logic functions per unit area. However, when scaled down to submicron size, the pass transistors will suffer from the subthreshold current problem. It is possible that when the fabrication technologies have been developed to enable scaling down to sub-micron devices, a technology such as complementary MOS, which does not draw any dc current, may be preferable to the nMOS technology used to illustrate this text. However, even if this occurs, the methodology developed in the text can still be applied in the design of integrated systems in that technology.

The limit to the kind of scaling described above occurs when the devices created are no longer able to perform the switching function. To perform the switching function, the ratio of transistor on to off conductance must be $> 1$, and therefore the voltage operating the circuit must be many times $kT/q$. For this reason, even circuits optimized for operation at the lowest possible supply voltages still require a VDD of ~ 0.5 volts. Devices in 1978 operate with a VDD of approximately five volts and minimum channel lengths of approximately six microns. Therefore, the kind of scaling we have envisioned here will take us to devices with approximately one half micron channel lengths and current densities approximately ten times what they are today. Power per unit area will remain constant over that range. Smaller devices might be built but must be used without lowering the voltage any further. Consequently the power per unit area will increase. Finally, there appears to be a fundamental limit of approximately one quarter micron channel length, where certain physical effects such as the tunneling through the gate oxide, and fluctuations in the positions of impurities in the depletion layers, begin to make the devices of smaller dimension unworkable.
References.


Reading References


Chapter 2: Integrated System Fabrication
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Sections:

- Patterning
- Scaling of Patterning Technology
- The Silicon Gate n-Channel Process
- Yield Statistics
- Scaling of the Processing Technology
- Design Rules
- Formal Description of Design Rules
- Electrical Parameters
- Current Limitations in Conductors
- A Closer Look at Some Details
- Choice of Technology

The series of steps by which a geometric pattern or set of geometric patterns is transformed into an operating integrated system is called a wafer fabrication process, or simply a process. An integrated system in MOS technology consists of a number of superimposed layers of conducting, insulating, and transistor forming materials. By arranging predetermined geometric shapes in each of these layers, a system of the required function may be constructed. The task of the integrated system designer is to devise the geometric shapes and their locations in each of the various layers of the system. The task of the process itself is to create the layers and transfer into each of them the geometric shapes determined by the system design.

Modern wafer fabrication is probably the most exacting production process ever developed. Since the 1950's, enormous human resources have been expended by the industry to perfect the myriad of details involved. The impurities in materials and chemical reagents are measured in parts per billion. Dimensions are controlled to a few parts per million. Each step has been carefully devised to produce some circuit feature with the minimum possible deviation from the ideal behavior. The results have been little short of spectacular: chips with many tens of thousands of transistors are being produced for under ten dollars each. In addition, wafer fabrication has reached a level of maturity where the system designer need not be concerned with the fine details of its execution. The following sections present a broad overview sufficient to convey the ideas involved, and in particular those relevant for system design. Our formulation of the basic concepts anticipates the evolution of the technology towards ever finer dimensions.

In this chapter we describe the patterning sequence and how it is applied in a simple, specific integrated system process: nMOS. A number of other topics are covered which are related to the processing technology, or are closely tied to the properties of the underlying materials.
Patterning

The overall fabrication process consists of the patterning of a particular sequence of successive layers. The patterning steps by which geometrical shapes are transferred into a layer of the final system, is very similar for each of the layers. The overall process is more easily visualized if we first describe the details of patterning one layer. We can then describe the particular sequence of layers used in the process to build up an integrated system, without repeating the details of patterning for each of the layers.

A common step in many processes is the creation of a silicon dioxide insulating layer on the surface of a silicon wafer, and the selective removal of sections of the insulating layer exposing the underlying silicon. We will use this step for our patterning example. The step begins with a bare polished silicon wafer, shown in cross section in figure 1. The wafer is exposed to oxygen in a high temperature furnace to grow a uniform layer of silicon dioxide on its surface, as shown in figure 2. After the wafer is cooled, it is coated with a thin film of organic resist material as shown in figure 3. The resist is thoroughly dried and baked to insure its integrity. The wafer is now ready to begin the patterning.

At the time of wafer fabrication the pattern to be transferred to the wafer surface exists as a mask. A mask is merely a transparent support material coated with a thin layer of opaque material. Certain portions of the opaque material are removed, leaving opaque material on the mask in the precise pattern required on the silicon surface. Such a mask with the desired pattern engraved upon it is brought face down into close proximity with the wafer surface, as shown in figure 4. The dark areas of opaque material on the surface of the mask are located where it is desired to leave silicon dioxide on the surface of the silicon. Openings in the mask correspond to areas where it is desired to remove silicon dioxide from the silicon surface. When the mask has been brought firmly into proximity with the wafer itself, its back surface is flooded with an intense source of ionizing radiation such as ultraviolet light or low energy x-rays. The radiation is stopped in areas where the mask has opaque material on its surface. Where there is no opaque material on the mask surface, the ionizing radiation passes on through into the resist, the silicon dioxide, and silicon. While the ionizing radiation has little effect on the silicon dioxide and silicon, it breaks down the molecular structure of the resist into considerably smaller molecules.

We have chosen to illustrate this text using positive resist, i.e. the resist material remaining after exposure and development corresponds to the opaque mask areas. Negative resists are also in
Figure 1. Bare Wafer

Figure 2. Oxidation

Figure 3. Coat w/Resist

Figure 4. Mask & Expose
Figure 5. Exposed Resist

Figure 6. Develop Resist

Figure 7. Etch

Figure 8. Remove Resist
common use. Positive resists are typically workable to finer feature sizes, and are likely to become dominant as the technology progresses.

After exposure to the ionizing radiation, the wafer has the characteristics shown in figure 5. In areas exposed to the radiation, the resist molecules have been broken down to much lighter molecular weight than that of unexposed resist molecules. The solubility of organic molecules in various organic solvents is a very steep function of the molecular weight of the molecules. Hence, it is possible to dissolve exposed resist material in solvents which will not dissolve the unexposed resist material. In this way the resist can be "developed" as shown in figure 6 by merely immersing the silicon wafer in a suitable solvent.

Thus far, the pattern originally existing as a set of opaque geometries on the mask surface has been transferred into a corresponding pattern in the resist material on the surface of the silicon dioxide. This same pattern can now be transferred to the silicon dioxide itself by exposing the wafer to a material which will etch silicon dioxide but not attack either the organic resist material or the silicon wafer surface. This etching step is usually done with hydrofluoric acid, which easily dissolves silicon dioxide. However, organic materials are very resistant to hydrofluoric acid, and it is incapable of etching the surface of silicon. The result of this etching step is shown in figure 7.

The final step in patterning is removal of the remaining organic resist material. Three techniques have been used to remove resist materials. Strong organic solvents will dissolve even unexposed resist material. Strong acids such as chromic acid actively attack organicics. The wafer can be exposed to atomic oxygen which will oxidize away any organic materials present on its surface. Once the resist material is removed, the finished pattern on the wafer surface is as shown in figure 8. Notice that we have transferred the geometric pattern which originally existed on the surface of the mask directly into the silicon dioxide on the wafer surface. While a foreign material was present on the wafer surface during the patterning process, it has now disappeared and the only materials present are those which will be part of the finished wafer.

A similar sequence of steps is used to selectively pattern each of the layers of the integrated system. These differ only in the details of the etchants used, etc. Thus as we study the processing of the various layers, the reader need not visualize all the details of the patterning sequence for each layer, but only recognize that a mask pattern for a layer can be transferred into a pattern in the material of that layer.
Scaling of Patterning Technology

As discussed in chapter 1, semiconductor devices could be at least an order of magnitude smaller in linear dimension than those typically manufactured in 1978 and still function correctly. The fundamental dimensional limitation is approximately a one quarter micron channel length, corresponding to a length unit \( \lambda \) (to be discussed under design rules) of approximately 0.1 micron. This limitation appears to apply to both bipolar and MOS technologies. It has been possible for several years to create sub-micron lines using electron beam and x-ray techniques, and there is considerable research and development under way to bring these patterning technologies into general manufacturing use. It appears that there are no fundamental barriers preventing creation of patterns for ultimately small devices. A more detailed discussion of the techniques involved is given in chapter 4.

The Silicon Gate n-Channel MOS Process

We now describe the particular sequence of patterned layers used to build up nMOS integrated circuits and systems. Figures 9 through 14 illustrate a simple but complete sequence of patterning and processing steps which are sufficient to fabricate a complete, integrated system. The example follows the fabrication of one simple circuit within a system, but all other circuits are simultaneously implemented by the same process. The example used is the basic inverter circuit. The top illustration in figures 9 through 14 shows the top view of the layers of the circuit layout. The lower illustration in each of those figures shows the cross section through the cut indicated by the downward arrows. The vertical scale in these cross sections has been greatly exaggerated for illustrative purposes.

The opening in the opaque material of the first mask is shown by the green outline in the top portion of figure 9. This opening exposes all areas that will eventually be the diffusion level. It includes the sources and drains of all transistors in the circuit, together with the transistor gate areas, and any diffusion level circuit interconnection paths. This mask is used for the first step in the process, the patterning of silicon dioxide on silicon as described in the previous section. The resulting cross section is shown in the lower portion of figure 9.

The second step in the process is to differentiate transistors which are normally "on" (depletion mode) from those which are normally "off" (enhancement mode). This is done by overcoating the wafer with resist material, exposing the resist material through openings in a second mask,
Fig. 9. Patterning SiO₂

Fig. 10. Patterning Ion Implantation

Fig. 11. Patterning Polysilicon

Fig. 12. Placing Diffused Region

Fig. 13. Placing Contact Cuts

Fig. 14. Patterning the Metal Layer
and developing it in the manner shown in figure 10. This patterning step leaves an opening in the resist material over the area to be selectively turned into depletion mode transistors. The actual conversion of the underlying silicon is then done by implanting ions of arsenic or antimony into the silicon surface. The resist material, where present, acts to prevent the ions from reaching the silicon surface. Therefore, ions are only implanted in the silicon area free of resist. The implanted layer, which causes a slight n-type conductivity in the underlying silicon, is shown by the yellow box in figure 10. Once the depletion areas are defined, the resist material is removed from the surface of the wafer.

The wafer is then heated while exposed to oxygen, to grow a very thin layer of silicon dioxide over its entire surface. It is then entirely coated with a thin layer of polycrystalline silicon, usually called polysilicon or poly for short. Note that this polysilicon layer is everywhere insulated from the underlying materials by the layer of thin oxide, and additionally by thicker oxide in some areas. The polysilicon will form the gates of all the transistors in the circuit and will also serve as a second layer for circuit interconnections. A third mask is used to pattern the polysilicon by steps similar to those previously described, with the result shown in red in figure 11. The left-most polysilicon area will function as the gate of the pull down transistor of the inverter we are constructing, while the square to the right will function as the gate of the depletion mode pull up transistor.

Once the polysilicon areas have been defined, n-type regions can be diffused into the p-type silicon substrate, forming the sources and drains of the transistors and the first level of interconnections. This step is done by first removing the thin gate oxide in all areas not covered by the (red) polysilicon. The wafer is then exposed to n-type impurities such as arsenic, antimony or phosphorus at high temperature for a sufficient period of time to allow these impurities to convert the exposed underlying silicon to n-type material. The areas of resulting n-type material are shown in green. Notice, in the cross section of figure 12, that the red polysilicon area and the thin oxide under it act to prevent impurities from diffusing into the underlying silicon. Therefore the impurities reach the silicon substrate only in areas not covered by the polysilicon and not overlain by the thick original oxide. In this way the active transistor area is formed in all places where the patterned polysilicon overlies the thin oxide area defined in the previous step. The diffusion level sources and drains of the transistors are automatically filled in between the polysilicon areas and extend up to the edges of the thick oxide areas. The major advantage of the silicon gate process is that it does not require a critical alignment between a mask which
defines the green source and drain areas and a separate mask which defines the gate areas. Rather, the transistors are formed by the intersection of the two masks, and the conducting n-type diffused regions are formed in all areas where the green mask is not covered by the red mask.

All the transistors of the basic inverter circuit are now defined. Connections must now be made to the input gate, between the gate and source of the pullup, and to VDD and GND. These interconnections will be made with a metal layer that can make contact to both the diffused areas and the polycrystalline areas. However, in order to ensure that the metal does not make contact to underlying areas except where intended, another layer of insulating oxide is coated over the entire circuit. At the places where the overlying metal is to make contact to either the polysilicon or the diffused areas, the overlying oxide is selectively removed by the patterning process as previously described. The result of coating the wafer with the overlying oxide and removing this oxide in places where contacts are desired, is shown in figure 13. In the top view, the black areas are those defined by openings in the contact mask, the fourth in the process's sequence of mask patterns. In cross section notice that in the contact areas all oxide has been removed down to either the polycrystalline silicon or the diffused area.

Once the overlying oxide has been patterned in this way, the entire wafer is coated with metal, usually aluminum, and the metal is patterned with a fifth mask to form the conducting areas required by the circuit. The top view in figure 14 shows three metal lines running vertically, the left most connecting to the input gate of the inverter, the center one being ground, and the right- most one forming the VDD connection to the inverter. The peculiar structure formed by the metal square slightly to the right of center connects the polysilicon gate of the depletion mode pull up transistor to its source and to the drain of the pull down transistor. Rather than making two separate contacts from the metal line to the pullup's polysilicon gate region and to the adjacent diffusion region, area can be conserved by coalescing the contacts into the compact arrangement shown. This geometrical arrangement is known as a butting contact and will be used extensively throughout the text.

In general, it is good practice to avoid placing contacts over active transistor area whenever possible. However, butting contacts in the location shown here reduce the area and simplify the geometry of the basic inverter and many other circuits, and have been so placed by the authors in many systems successfully implemented by a number of different commercial wafer fabrication lines. A more conservative approach would be to place the butting contact adjacent to, rather than over, the active pullup area. See also the later section on design rules in this chapter.
The inherent properties of the silicon gate process allow the blue metal layer to cross over either the red polysilicon layer or the green diffused areas, without making contact unless one is specifically provided. The red polysilicon areas, however, cannot cross the green diffused areas without forming a transistor. The transistors formed by the intersection of these two masks can be either enhancement mode if no yellow implantation is provided, or depletion mode if such an implantation is provided. Hence, the enhancement mode transistors are defined by the intersection of the green and red masks while the depletion mode transistors are defined by the intersection of the green, red and yellow masks.

If we wish to fabricate only a small number of prototype system chips and to have access to the metal level for the probing of test points, the wafer fabrication sequence can be terminated at this step. However, when fabricating large numbers of chips of a debugged design, the wafer surface is usually coated with another layer of oxide. This step, called overglassing, provides physical protection for the devices in the system. A sixth mask is then used to pattern contact cuts in the overglassing at the locations of relatively large metal contact pads.

Each wafer contains many individual chips. The chips are separated by scribing the wafer surface with a diamond scribe, and then fracturing the wafer along the scribe lines. Each individual chip is then cemented in place in a package, and fine metal wire leads are bonded to the metal contact pads on the chip and to pads in the package which connect with its external pins. A cover is then cemented over the recess in the package which contains the silicon chip, and the completed system is ready for testing and use.

Yield Statistics

Of the large number of individual integrated system chips fabricated on a single silicon wafer, only a fraction will be completely functional. Flaws in the masks, dust particles on the wafer surface, defects in the underlying silicon, etc., all cause certain devices to be less than perfect. With present design techniques, any single flaw of sufficient size will kill an entire chip.

The simplest model for the yield, or the fraction of the chips fabricated which do not contain fatal flaws, assumes (naively) that the flaws are randomly distributed over the wafer, and that one or more flaws anywhere on a chip will cause it to be non-operative. If there are N fatal flaws per unit area, and the area of an individual chip is A, the probability that a chip has n flaws is in the
simplest case just given by the Poisson distribution, \( P_n(NA) \). The probability of a good chip is:

\[
P_n(NA) = e^{-NA}
\]

While this equation does not accurately represent the detailed behavior of real fabrication processes, it is a good approximate model for estimating the yield of alternative designs. The exponential is such a steep function that a very simple rule is possible: chips with areas many times \( 1/N \) will simply never be found without flaws. Areas must be kept less than a few times \( 1/N \) if one flaw will kill a system. Design forms may be developed in the future which will permit systems to work even in the presence of flaws. If such forms are developed, the entire notion of yield will be completely changed and much larger chips will be possible.

Once a wafer has been fabricated, each chip must be tested to determine if it is functional. Testing of simple combinatorial logic networks is straightforward and may be done completely. Complete testing of complex systems with internal sequencing is not in general possible, and most integrated system chips manufactured, even at 1978 levels of complexity, are not economically testable even for a small fraction of their possible internal states.

As time passes and the number of devices per chip increases, it will become important to consider including special functions in the design of integrated systems to improve their testability. The basic problem is to linearize an otherwise combinatorial problem. One approach to this is:

(i) Define the entire system as a set of register to register transfer blocks, i.e. successive stages of storage registers with combinational logic between them.

(ii) Provide for reading and writing from the external world to/from each of the storage registers.

The storage locations are first tested independently for their ability to store data or control information. If all storage locations pass this test, each combinational logic block can be tested separately, by use of its input and output storage locations. Such a test becomes essentially linear in the number of components, and may be accomplished in an acceptable time period, even for extremely complex systems. However, without access to the individual storage locations, testing rapidly becomes hopeless. For this reason even present day microprocessors are very incompletely tested. When one is used for a while, an apparently new and sudden malfunction may simply be the first occurrence of a particular state of control and data in the system, and thus may represent the first time the device had been "tested" under those conditions.
From experience gained in testing memory parts, it is known that the behavior of one circuit can be influenced by the state of a nearby circuit. For example, a memory cell may be able to remember both a logic-1 and a logic-0 if its neighbor is at a logic-0, but may be able to retain only a logic-0 if its neighbor is at a logic-1. Failures of this type are dependent upon the data patterns present in the system, and are known as pattern sensitive failures. In a reasonable (or even an unreasonable) time, it is not possible to exercise even a minute fraction of all the combinations of bit patterns of many integrated systems. What is done instead is to apply our knowledge of the physics of such failures, and construct a model for possible failure modes. In the memory example, we may conclude that any flaw not visible optically will be unable to reach beyond the immediate locality of the cell involved. Hence, pattern sensitivity in the behavior of a particular cell may be introduced by other cells in the same row or column of an array of memory cells, or by diagonal nearest neighbors. A test for pattern sensitivity under this model is quite fast, being only slightly worse than a linear function of the number of devices on the chip.

In order to test for pattern sensitive failures, we must construct a physical model for the possible failure mechanisms. This model will inevitably include the physical proximity of other signals. For this reason, any practical test for pattern sensitive failures must be based on a knowledge of the physical location of the various elements of the subsystem being tested. The task of preparing such tests is thus greatly eased by regularity in the design and physical layout of a system.

Scaling of the Processing Technology

In order to have a complete process for sub-micron transistors, it is necessary not only to make patterns in the resist material but to transfer these patterns to the underlying layers in the silicon and silicon dioxide. Traditionally, wet etching processes have been used. However, wet etching processes do not scale well into the sub-micron range.

Alternatives are currently being developed which appear workable. Etching with plasmas (i.e. glow discharges of gaseous materials resulting in free ions of great chemical activity) is already used in a number of advanced processing facilities. It is known that very well controlled etching can be achieved in this way and it seems likely that essentially no wet processing will be used in the construction of sub-micron devices. Ion implantation, an ideal method for achieving controlled doses of impurity ions in the silicon surface, is already a common production technique in essentially all MOS processing facilities.
Metal layers for sub-micron processes must be thicker in relationship to their width than today's commercial processing technology allows. A possible solution to this problem may be the use of a process known as ion milling for metal patterning. In this process, ions of modest energy sputter away any metal not covered with resist material, yielding much steeper sides on the metal thus patterned than do current wet etching processes.

It appears that the basic technological pieces exist to enable development of a complete patterning and wafer fabrication process at sub-micron dimensions. In reality, the ultimate submicron process will not emerge full-blown, but dimensions will gradually be reduced, as one after another of the myriad of technological difficulties are surmounted. The sketch we have given is rather an artist's conception of the possibility of such an ultimate process. We do believe, however, that the evolution of this process is of fundamental importance to the entire electronics industry.
Design Rules

Perhaps the most powerful attribute of modern wafer fabrication processes is that they are pattern independent. That is, there is a clean separation between the processing done during wafer fabrication and the design effort which creates the patterns to be implemented. This separation requires a precise definition to the designer of the capabilities of the processing line. This specification usually takes the form of a set of permissible geometries which may be used by the designer with the knowledge that they are within the resolution of the process itself and that they do not violate the device physics required for proper operation of transistors and interconnections formed by the process. When reduced to their simplest form, such geometrical constraints are called design rules. The constraints are of the form of minimum allowable values for certain widths, separations, extensions, and overlaps of geometrical objects patterned in the various levels of a system.

As processes have improved over the years, the absolute values of the permissible sizes and spacings of various layers have become progressively smaller. There is no evidence that this trend is abating. In fact, there is every reason to believe that at least another order of magnitude of shrinkage in linear dimensions is possible. For this reason we present a set of design rules in dimensionless form, as constraints on the allowable ratios of certain distances to a basic length unit. The basic unit of length measurement used is equal to the fundamental resolution of the process itself. This is the distance by which a geometrical feature on any one layer may stray from another geometrical feature on the same or on another layer, all processing factors considered and an appropriate safety factor added. It is set by phenomena such as overetching, misalignment between mask levels, distortion of the silicon wafer ("runout") due to high temperature processing, over or underexposure of resist, etc. All dimensions are given in terms of this elementary distance unit, which we call the length-unit, \( \lambda \). In 1978 the length-unit \( \lambda \) is approximately 3 microns for typical commercial processes. One micron (\( \mu m \)) = \( 10^{-6} \) meters.

The rules given below have been abstracted from a number of processes over a range of values of \( \lambda \), corresponding to different points in time at different fabrication areas. They represent somewhat of a "least common denominator" likely to be representative of nMOS design rules for a reasonable period of time, as the value of \( \lambda \) decreases in the future.

A typical minimum for the line width \( W_d \) of the diffused regions is \( 2\lambda \), as shown in figure 15. The spacing required between two electrically separate diffused regions is a parameter which
depends not merely upon the geometric resolution of the process, but also upon the physics of the devices formed. If two diffused regions pass too close together, the depletion layers associated with the junctions formed by these regions may overlap and result in a current flowing between the two regions when none was intended. In typical processes a safe rule of thumb is to allow $3\lambda$ of separation, $S_{dd}$, between any two diffused regions which are unconnected, as shown in figure 16. The width of a depletion layer associated with any diffused region depends upon the voltage on the region. If one of the regions is at ground potential, its depletion layer will of necessity be quite thin. In addition some processes provide a heavier doping level at the surface of the wafer between the diffused areas in order to alleviate the problem of overlap of depletion layers. In cases where either very low voltage exists on both diffused regions or where a heavily doped region has been implanted in the surface between the diffused areas, it is often possible to space diffused areas $2\lambda$ apart. However, this should not be done without carefully checking the actual process by which the design is to be fabricated.

The minimum for the width $W_p$ of polysilicon lines is similarly $2\lambda$. No depletion layers are associated with polysilicon lines, and therefore the separation, $S_{pp}$, of two such lines may be as little as $2\lambda$. These rules are illustrated in figures 17 and 18.

We have so far considered the diffused and polysilicon layers separately. Another type of design rule concerns how the two layers interact with each other. Figure 19 shows a situation where a diffused line is running parallel to an independent polysilicon line, to which it is not anywhere connected. The only consideration here is that the two unconnected lines not overlap. If they did they would form an unwanted capacitor. Avoidance of this overlap requires a separation $S_{pd}$ of only $\lambda$ between the two regions as shown in figure 19. A slightly more complex situation is shown in figure 20, where a polysilicon gate area intentionally crosses a diffused area, thereby forming a transistor. In order to make absolutely sure that the diffused region does not reach around the end of the gate and short out the drain to source path of the transistor with a thin diffused area, it is necessary for the polysilicon gate to extend a distance $E_{pd}$ of at least $2\lambda$ beyond the nominal boundary of the diffused area, as shown in figure 20.

A composite of several of these design rules is shown in figure 21. Note that the minimum width for a diffused region applies to diffused regions formed between a normal boundary of the diffused region and an edge of a transistor as well as to a diffused line formed by two normal boundaries. This situation is illustrated in the lower left corner of the figure.
Fig. 15. $W_d/\lambda \geq 2$

Fig. 16. $S_{dd}/\lambda \geq 3$

Fig. 17. $W_p/\lambda \geq 2$

Fig. 18. $S_{pp}/\lambda \geq 2$

Fig. 19. $S_{pd}/\lambda \geq 1$

Fig. 20. $E_{pd}/\lambda \geq 2$

Fig. 21. Example of Several Rules

Fig. 22. $S_{ig}/\lambda \geq 1/2$  $E_{ig}/\lambda \geq 1/2$
Fig. 23. $W_c / \lambda \geq 2$, $F_{dc} / \lambda \geq 1$

Fig. 24. $F_{pc} / \lambda \geq 1$

Fig. 25. $S_{cc} / \lambda \geq 2$, $S_{cg} / \lambda \geq 2$

Fig. 26. $O_{pd} / \lambda = 1$, and details of butting contact

Fig. 27. $W_m / \lambda \geq 3$, $S_{mm} / \lambda \geq 3$

Fig. 28. $F_{mc} / \lambda \geq 1$
As we have seen in figure 10, ion implantation in the region which becomes the gate of a transistor will convert the resulting transistor into the depletion mode type. It is important that the implanted region extend outward beyond all four boundaries of the gate region, as shown in figure 22. To avoid any possibility that some small fraction of the transistor might remain enhancement mode, the yellow ion implantation region should extend a distance $E_{ig}$ of at least $1\frac{1}{2} \lambda$ beyond each edge of the gate region. The separation $S_{ig}$ between an ion implantation region and an adjacent enhancement mode transistor gate region should also be at least $1\frac{1}{2} \lambda$. Both situations and their design rules are illustrated in the figure.

A contact may be formed between the metal layer and either the diffused level or the polysilicon level by means of the contact mask. A set of rules apply to the amount by which each layer must provide an area surrounding any contact to it, so that the contact opening not find its way around the layer to something unintended below it. Since no physical factors apply here other than the relative registration of two levels, a very simple set of design rules results. Each level involved in a given contact must extend beyond the outer boundary of the contact cut by $\lambda$ at all points, as illustrated by extension distances $E_{dc}$, $E_{pc}$, and $E_{mc}$, in figures 23, 24, and 26. The contacts themselves, like the minimum width lines in the other levels, must be at least $2\lambda$ long and $2\lambda$ wide ($W_c$). This situation is illustrated for the diffusion and polysilicon levels in figures 23 and 24. When making contact between a large metal region and a large diffused region, many small contacts spaced $2 \lambda$ apart should be used, as shown in figure 25. Contact cuts to diffusion should be at least $2\lambda$ from the nearest gate region, as shown in figure 25.

Note that a cut down to the polysilicon level does not penetrate the polysilicon. Thus one can in principle make a contact cut to poly over a gate region, and such contacts are permitted in these design rules. However, since such a cut must be $2\lambda$ wide and surrounded on all sides by $1\lambda$ of poly, it is not possible to make such a contact above a minimum size transistor's gate region. Also, as device dimensions scale down and the poly and thin oxide become ever thinner, such cuts might penetrate too far, and thus they may not be allowed in the design rules in the future.

When a direct connection is required between a polysilicon region and a diffused region, we normally use a construct known as the butting contact. The detailed geometric layout of the butting contact is shown in figure 26. In its minimum sized configuration, it is composed of a square region of diffision $4 \lambda$ on a side, overlapped by a $3 \lambda$ by $4 \lambda$ rectangle of polysilicon. A rectangular contact cut, $2 \lambda$ by $4 \lambda$ in size, is made in the center of this structure. The structure is then overlayed with metal, thus connecting the polysilicon to the diffusion. The rules involved
in figure 26 are identical to those given so far, with the addition of a minimum of one \( \lambda \) overlap, \( \Omega_{Dd} \), of the diffused and polysilicon layers in the center area of the contact.

In considering the design rules for the metal layer, notice that this layer in general runs over much more rugged terrain than any other level, as can be seen by referring to the cross section of figure 14. For this reason it is generally accepted practice to allow somewhat wider minimum lines and spaces for the metal layer than for the other layers. As a good working rule 3\( \lambda \) widths (\( W_m \)), and 3\( \lambda \) separations (\( S_{mm} \)) between independent metal lines should be provided, as shown in figure 27.

The metal layer must surround the contact layer in much the same way that the diffused and polysilicon layers did. Since the resist material used for patterning the metal generally accumulates in the low areas of the wafer, it tends to be thicker in the neighborhood of contact than elsewhere. For this reason metal tends to be slightly larger after patterning in the vicinity of a contact than elsewhere. It is generally sufficient to allow only one \( \lambda \) of space around the contact region for the metal, as for the other two layers. The rule for metal surrounding contacts is shown in figure 28.

Additional layout artifacts, and guidelines and rules related to the layout artifacts, such as alignment marks, which are associated with conveying a chip's layout through the processes of maskmaking and wafer fabrication are given in chapter 4. Included there are guidelines for sizing such macroscopic layout artifacts as chip scribe lines, wire bonding pads, etc. However, the design rules given here in chapter 2 are sufficient for the layout of the functional circuitry within an nMOS integrated system.

The above design rules are likely to remain valid as the length-unit \( \lambda \) scales down in size with the passage of time. Occasionally, for specific commercial fabrication processes, some one or more of these rules may be relaxed or replaced by more complex rules, enabling slight reductions in the area of a system. While these details may be important for certain competitive products such as memory systems, they have the disadvantage of making the system design a captive of the process specific design rules. Extensive redesign and checking is required to scale down such a design as the length-unit scales down. For this reason, we recommend use of the dimensionless rules given, especially for prototype systems. Designs implemented according to these rules are easily scaled, and may have reasonable longevity.
Formal Description of Design Rules

{ in preparation }

Electrical Parameters

By satisfying the constraints imposed by the design rules, designers may create circuit layout patterns with the knowledge that the appropriate transistors, lines, etc., produced by the wafer fabrication process will be as originally specified in their layout patterns. To complete a design it is necessary to also know the electrical parameters of the transistors, diffused layers, polysilicon layers, etc., so that the performance of circuits can be evaluated. The resistances per square of the various layers and the capacitance per square micron with respect to underlying substrate are shown in Table 1. Note that the resistance of a square of material contacted along two opposite sides is independent of the size of the square, and equals the resistivity of the material divided by its thickness. The tabulated values are typical of processes running in 1978. As the circuit dimensions are scaled down by dividing by a factor \( \alpha \), the parameters scale approximately as shown in the table.

<table>
<thead>
<tr>
<th>Resistances:</th>
<th>Metal</th>
<th>\sim 0.1 \text{ohms/\square}</th>
<th>Resistances/square scale up by ( \alpha ).</th>
</tr>
</thead>
<tbody>
<tr>
<td>Diffusion</td>
<td>\sim 10 \text{ohms/\square}</td>
<td>as dimensions scale down by ( \alpha ).</td>
<td></td>
</tr>
<tr>
<td>Poly</td>
<td>\sim 15-100 \text{ohms/\square}</td>
<td>except that the transistor</td>
<td></td>
</tr>
<tr>
<td>Transistor</td>
<td>\sim 10^4 \text{ohms/\square} R/\square</td>
<td>is independent of ( \alpha ).</td>
<td></td>
</tr>
</tbody>
</table>

<table>
<thead>
<tr>
<th>Capacitances:</th>
<th>Gate-channel</th>
<th>\sim 4 \times 10^{-4} \text{pf/\mu m}</th>
<th>Capacitances/micron$^2$ scale</th>
</tr>
</thead>
<tbody>
<tr>
<td>Diffusion</td>
<td>\sim 0.8 \times 10^{-4} \text{pf/\mu m}</td>
<td>up by ( \alpha ), as dimensions</td>
<td></td>
</tr>
<tr>
<td>Poly</td>
<td>\sim 0.4 \times 10^{-4} \text{pf/\mu m}</td>
<td>scale down by ( \alpha )</td>
<td></td>
</tr>
<tr>
<td>Metal</td>
<td>\sim 0.3 \times 10^{-4} \text{pf/\mu m}</td>
<td></td>
<td></td>
</tr>
</tbody>
</table>

Table 1. Typical MOS Electrical Parameters (1978).

The relative resistance values of metal, diffusion, poly, and drain to source paths of transistors are quite different. Diffusion and good polysilicon layers have approximately one hundred times the
resistance per square area of the metal layer. A fully turned on transistor has approximately one thousand times the resistance of the diffused and polysilicon layers. The capacitances are not as wildly different as the resistances of the various layers. Compare the capacitances in Table 1 to the gate to channel capacitance, as a reference. The diffused areas typically have one fifth the capacitance per square micron. Polysilicon on thick oxide has approximately one tenth, and the metal layer slightly less than one tenth, of the gate-channel capacitance per square micron.

The relative values of the resistances and capacitances are not expected to vary dramatically as the processes evolve towards smaller dimensions, with the exception of the transistor resistance per square, which is independent of $\alpha$.

One note of warning: There is a wide range of possible values of polysilicon resistance for different commercial processes. Polycrystalline silicon suffers from inordinately high resistances at the crystal grain boundaries if the doping level in the polysilicon itself is not held quite high. This disease does not affect the diffused layers. For this reason, any processing which tends to degrade the doping levels in the diffused and polysilicon layers, affects the polysilicon resistance much more dramatically than the resistance of the diffused area. It is in general difficult to design circuits which are optimum over the entire range of polysilicon resistivity. If a circuit is to be run on a variety of fabrication lines, it is desirable for the circuit to be designed in such a way that no appreciable current is drawn through a long thin line of polysilicon. In an important example in Chapter 5., polysilicon lines are used as buses along which information flows. The timing of these buses can be dramatically affected by the resistance of the polysilicon. However, the protocol used on these busses has the polysilicon lines precharged during one period of a clock and then pulled low by the appropriate bus source during a following clock period. In this way the circuit is guaranteed to work independent of the resistance of the poly. However, it may be considerably slower in processes of high poly resistivity.

Current Limitations in Conductors

One limit which is not covered in either the design rules or the electrical parameters section is that associated with the maximum currents through metal conductors. There is a physical process called metal migration whereby a current flux through a metal conductor, exceeding a certain limit, causes the metal atoms to move slowly in the direction of the current. If there is a small
constriction in the metal, the current density will be higher and therefore more metal atoms will be carried forward from that point, narrowing the point still more. Hence, metal migration is a destructive mechanism causing open circuits in the metal layer carrying heavy currents.

For metals like aluminum this limit is a few times $10^5$ amperes per square centimeter, i.e. a few milliamperes per square micron of cross section. This limit does not interfere too drastically with the design of integrated systems in current MOS technologies. However, many metal conductors in present integrated systems are operated near their current limit, and currents do not scale well as the individual elements are made smaller. Applying the scaling rules developed earlier, we found that the power per unit area is independent of the scale down ratio. However, the supply voltage decreases and therefore the current per unit area increases as the devices are scaled down. For this reason it will not be possible to use processes for very large scale integrated systems where the metal thickness scales in the same way as other dimensions in the circuit. Much work will likely be done to develop processes enabling fabrication of metal lines of greater height relative to width than is presently possible.

Short pulses of current are known to contribute much less to metal migration than steady direct current. Nanosecond pulses of currents two orders of magnitude higher than the dc limit given above may be carried in metal conductors without apparent damage. Therefore, switching current may not be as damaging to metal conductors as a steady current.

These effects strongly favor processes like CMOS which do not require static dc current, and favor design methodologies which maximize system function per unit dc current.

A Closer Look at Some Details

Thus far our discussion of fabrication has been a general one, adequate for readers whose primary interest is in the systems aspects of VLSI. The following sections involve a more detailed examination of the capacitance of several important structures and a discussion of the relative merits and scaling behavior of several common processes. We suggest that the reader just skim through these sections during the first reading of this text.

In Table 1 we gave typical capacitance for the various layers to the substrate. These capacitances are those which would be measured if the voltage on the particular layer were zero (relative to the substrate). The dependence upon voltage of the capacitance of the different layers may
sometimes be important and we will now discuss how this dependence arises. References R1, R2, R4, and Reference 4 of Chapter 1, are good sources for those wishing more background information on the concepts of device physics used in this text.

When a negative voltage is applied to an n-type diffused region relative to the p-type bulk silicon, the negative electrons are pushed out of the n-type layer into the bulk and a current flows. In integrated systems we are careful to never allow the voltage on the n-type diffused regions to be more negative than the p-type bulk. Diffused regions are biased positively with respect to the p-type bulk, resulting in a reversed biased p/n junction. With the exception of a small leakage current, the reverse biased p/n junction acts merely to isolate one diffused region from another. The p-type bulk of our integrated system has a small number (typically $10^{15}$-$10^{16}$ per cubic centimeter) of impurity atoms. When a voltage is applied to an n-type diffused region, its influence is felt well out into the p-type bulk. Positive charge carriers in the p-type bulk are repelled from the positively charged n-type layer, thereby exposing negatively charged impurity ions. The region surrounding the n-type diffused layer which has been depleted of positive charge carriers is referred to as a depletion layer and is shown schematically in Figure 29b. As the voltage on the n-type layer is increased, charge carriers are pushed further back from the junction between the n-type layer into the p-type bulk, widening the depletion layer and exposing more charged impurity ions. The charge thus induced in the depletion layer as the voltage on the n-type diffused region is increased is responsible for the capacitance of the n-type diffused region relative to the substrate.

We will now consider a unit area of the junction. The total charge in the depletion layer per unit area is proportional to the number per unit volume of impurity ions in the bulk ($N$), and the width, $s_0$, of the depletion layer.

$$\text{Total charge/area } \propto Ns_0$$

The electric field in the region is proportional to the charge per unit area.

$$\text{Electric field } \propto \text{charge/area } \propto Ns_0$$

The voltage between the n-type diffused layer and the p-type bulk on the far side of the depletion layer is proportional to the electric field times thickness of the depletion layer, and
Fig. 29a. n-type Diffusion in p-type Bulk Silicon

Fig. 29b. Depletion Layer

Fig. 30. C/A as fcn(V)

Fig. 31. C/A as fcn(V_{applied})

Fig. 32. Capacitance of Poly or Metal over Thick Oxide

Fig. 33. MOS Gate Capacitance as fcn(V)

(graphs, press)
therefore to the density of negatively charged ions in the depletion layer times the square of the width of the depletion layer.

\[ \text{Voltage } \propto \text{electric field } \times s_0 \propto Ns_0^2 \]

The capacitance per unit area is just the charge per unit area divided by the voltage across the depletion layer. From the above equations the capacitance is proportional to the square root of the density of impurity atoms in the p-type bulk divided by the voltage.

\[ \text{Capacitance/area } = \frac{Q}{V} \propto \frac{1}{s_0} \propto \left(\frac{N}{V}\right)^{1/2} \]

This relationship is plotted in Figure 30. Notice that the capacitance tends towards infinity as the voltage across the junction tends to zero. It would seem that this large capacitance would be disastrous for the performance of our integrated systems. However, this proves not to be the case. When the p/n junction was formed the n-type region had an excess of negative charge carriers while the p-type bulk had an excess of positive charge carriers. When the two were brought together to form the junction, there was no voltage to prevent charge carriers of either type from flowing over into the opposite region. This initial flow caused the n-type layer to become more positive than the p-type layer. This flow ceased when just enough voltage built up to stop it. In silicon the voltage required to prevent the flow of charge carriers in such a situation is approximately 0.7 volt. Thus the true voltage across the junction is this initial "built-in" voltage plus the voltage we apply in our circuit. The variation of the capacitance per unit area with applied voltage is shown in Figure 31. An approximate equation which can be used to calculate the junction capacitance \( C_j \) per unit area of diffused layers as a function of the applied voltage is given by:

\[ C_j = 4.5 \times 10^{-12} [N/(V + 0.7)]^{1/2} \text{ pF/\mu m}^2 \]

In this equation, \( N \), the density of impurity ions in the p-type bulk, should be given in number per cm\(^3\). The voltage is in volts and the capacitance per unit area is evaluated in picofarads per square micron. This equation is adequate for most design purposes.

Aside from the diffused regions, there are two other situations where the capacitance is of interest. The first is poly or metal over thick oxide and the second is the gate of an MOS
transistor. We will discuss poly or metal over oxide first. Figure 32 illustrates once more the capacitance per unit area of a junction over the p-type bulk. If the poly or metal layer was lain on an oxide much thinner than the depletion layer, its capacitance would be nearly the same as that of the corresponding p/n junction. However, if an oxide is interposed whose thickness is of the order of the depletion layer thickness, the capacitance of the poly or metal line will be decreased. The formula which applies in this case is given by:

\[
\frac{1}{C_{\text{total}}} = \frac{1}{C_j} + \frac{1}{C_{\text{ox}}}
\]

A typical dependence is shown in Figure 32. For an oxide thickness \(d\), \(C_{\text{ox}} = 3.5 \times 10^{-2}/d\), where the thickness \(d\) is given in angstrom units \(10^{-4}\) microns, and the result is in picofarads per square micron as before.

The most spectacular voltage dependence of a capacitance in the technology we will be using is that of the gate of an MOS transistor. When the gate voltage \(V_{gs}\) is less than the threshold voltage \(V_{th}\), the capacitance of the gate to the bulk is just that given above for metal or poly over oxide, since the voltage on the gate merely depletes positive charge carriers back from the channel area. However, when the voltage on the gate reaches the threshold voltage of the transistor, negative charge is brought in under the gate oxide from the source of the transistor and the capacitance changes abruptly from the small value associated with depleting charges in the bulk to the much larger oxide capacitance between the gate and the channel region. Further increase in voltage on the gate merely increases the amount of mobile charge under the gate oxide with no change in the width of the depletion layer underneath the channel. Hence, the character of the gate capacitance changes abruptly as the gate voltage passes through the threshold voltage.

The dependence of the total gate capacitance on gate voltage is shown in figure 33. The capacitance from channel to bulk is completely separate from the gate to channel capacitance. It is associated with the depletion layer underneath the channel region, and is almost identical to that of a diffused region of the same area. When the gate voltage is below threshold, the gate to channel capacitance disappears altogether leaving only the small parasitic overlap capacitances between the gate and the source and drain regions.
Choice of Technology

Before proceeding to the chapters on system design, let us briefly examine some alternative technologies. Using the knowledge developed in these first two chapters, we will discuss the reasons for selecting nMOS as the single technology used to illustrate integrated systems in this text. Some of the factors which must be considered in choosing a technology include circuit density, richness of available circuit functions, performance per unit power, the topological properties of circuit interconnection paths, suitability for total system implementation, and general availability of processing facilities.

As the technology advances, more system modules can be placed on the same sized chip. An ultimate goal is the fabrication of large scale systems on single chips of silicon. For this goal to be attained, any signal which is required in the system other than inputs, outputs, VDD, and GND, must be generated in the technology on the chip. In other words, no subsystem can require a different technology for the generation of its internal signals. Thus such technologies as magnetic bubbles are ruled out for full integrated systems because they are not able to create the signals required for all operations in the on-chip medium.

We believe that for any silicon technology to implement practical large scale systems, it must provide two kinds of transistor. The rationale for this observation is as follows. In order to provide some kind of nonlinear threshold phenomenon there must be a transistor which is normally off when its control input is at the lowest voltage used in the system. Bipolar technologies use NPN transistors for this purpose. The nMOS technology uses n-channel enhancement mode devices. In addition to this transistor, a separate type of transistor must be supplied to allow the output of a driver device to reach the highest voltage in the circuit (VDD). In the bipolar technologies, PNP lateral devices are used to supply this function, in the n-channel technology a depletion mode device is used, and in complementary MOS technology a p-channel enhancement mode device is used. All three choices allow output voltages of drivers to reach VDD and thus meet the above criterion.

To date three technologies have emerged which are reasonably high in density and scale to submicron dimensions without an explosion in the power per unit area required for their operation. These are the n-channel silicon gate process, the complementary MOS silicon gate process, and the integrated injection logic, or I$^2$L, process. Although present forms of I$^2$L technology lack the additional level of interconnect available in the silicon gate technologies, there
is no inherent reason that such a level could not be provided. It is important to note that increasing the flexibility of interconnect enriches the types of array functions which can be created. $I^2L$ has the advantage over nMOS that the power per unit area (and hence the effective $\tau$ of its elementary logic functions) can be controlled by an off-chip voltage. The decision concerning what point on the speed vs power curve to operate may thus be postponed until the time of application (or even changed dynamically).

The nMOS scaling has been described previously. Any technology in which a capacitive layer on the surface induces a charge in transit under it to form the current control "transistor" will scale in the same way. Examples include Schottky Barrier Gate FET's (MESFETs), Junction FET's, and CMOS.

There are certain MOS processes (VMOS, DMOS) of an intermediate form in which the channel length is determined by diffusion profiles. While competitive at present feature sizes, these are likely to be interim technologies which will present no particular advantage at submicron feature sizes.

Scaling of the bipolar technology\(^1\) is quite different from that of MOS technologies. For completeness, we include here a discussion of the scaling of bipolar devices, which may be of interest to those familiar with those technologies.

Traditionally, bipolar circuits have been "fast" because their transit time was determined by the narrow base width of the bipolar devices. In the 1950's, technologists learned how to form bipolar transistor base regions as the difference between two impurity diffusion profiles. This technique allowed very precise control of the distance perpendicular to the silicon surface, and therefore permitted the construction of very thin base regions with correspondingly short transit times. Since current in a bipolar device flows perpendicular to the surface, both the current and the capacitance of such devices are decreased by the same factor as the device surface dimensions are scaled down, resulting in no change in time performance. The base widths of high performance bipolar devices are already nearly as thin as device physics allows. For this reason, the delay times of bipolar circuits is expected to remain approximately constant as their surface dimensions are scaled down.

The properties of bipolar devices may be analyzed as follows. The collector current is due to the diffusion of electrons from emitter to collector. For a minority carrier density $N(x)$ varying
linearly with distance \( x \), from \( N_0 \) at the emitter to zero at the collector (at \( x = d \)), the current \( I \) per unit area \( A \) is:

\[
\frac{I}{A} = q(2D)d\frac{N}{dx} = q(2D)N_0/d = q(2kT/q)\mu N_0/d \quad \text{[eq.2]}
\]

where the diffusion constant \( D = \mu kT/q \). The factor of two multiplies the diffusion constant in eq.2 because high performance bipolar devices operate at high injection level (injected minority carrier density much greater than equilibrium majority carrier density). The inherent stored charge in the base region is:

\[
Q/A = N_0 d/2 \quad \text{[eq.3]}
\]

Therefore, the transit time is:

\[
\tau = Q/I = d^2/[4\mu kT/q] \quad \text{[eq.4]}
\]

The form of equation 4 is exactly the same as that for MOS devices (eq.1., chapter 1.), with the voltage in the bipolar case being equal to \( 4kT/q \) (at room temperature \( kT/q = 0.025 \) volts). A direct comparison of the transit times is shown in Table 2.

<table>
<thead>
<tr>
<th>Table 2. Transit Time:</th>
</tr>
</thead>
<tbody>
<tr>
<td>( \tau = \frac{(\text{Distance})^2}{(\text{Mobility} \times \text{Voltage})} )</td>
</tr>
</tbody>
</table>

<table>
<thead>
<tr>
<th>MOSFET:</th>
<th>MESFET, JFET:</th>
<th>Bipolar:</th>
</tr>
</thead>
<tbody>
<tr>
<td>\text{Distance:}</td>
<td>channel length</td>
<td>channel length</td>
</tr>
<tr>
<td>\text{Voltage:}</td>
<td>\text{~ VDD/2},</td>
<td>\text{~ VDD/2},</td>
</tr>
<tr>
<td></td>
<td>many kT/q</td>
<td>many kT/q</td>
</tr>
<tr>
<td>\text{Mobility:}</td>
<td>surface mobility.</td>
<td>bulk mobility.</td>
</tr>
<tr>
<td>( \text{cm}^2/\text{v-sec}, \text{(Si)} )</td>
<td>( \sim 800 )</td>
<td>( \sim 1300 )</td>
</tr>
</tbody>
</table>

At the smallest dimensions to which devices can be scaled, the base width of bipolar devices and the channel length of FET devices are limited by the same basic set of physical constraints, and are therefore similar in dimension. The voltage on the FET devices must be many times \( kT/q \) to achieve the required nonlinearity. Hence at ultimately limiting small dimensions the two types of
device have roughly equivalent transit times. At these limiting dimensions, choices between competing technologies will be made primarily on the grounds of the topological properties of their interconnects, the functional richness of their basic circuits, simplicity of process, and ability to control dc current per unit area. As supply voltages are scaled down to the 1 volt range, MOS devices become similar in most respects to other FET type devices, and it is possible that mixed forms (MOS-JFET, MOS-MESFET, Bipolar-MESFET, etc.) may emerge as the ultimate integrated system technologies.

We have chosen to illustrate this text with examples drawn from the n-channel silicon gate depletion mode load technology. The reasons for this choice in 1978 are quite clear. In addition to meeting the required technical criteria we have described, this technology provides some important practical advantages to the student and to the teacher. It is the only high density technology which has achieved universal acceptance across company and product boundaries. Readers wishing to implement integrated system designs may have wafers fabricated by essentially any wafer fabrication firm, without fear that slight changes in the process or the vagaries of relationships with a particular firm will cut off their source of supply. It is also presently the highest density process available. This certainty of access to fabrication lines, the more generally widespread knowledge of nMOS technology among members of the technical community, its density, and its performance similarity with bipolar technology in its ultimate scaling, are all important factors supporting its choice for this text on VLSI Systems. However, the principles and techniques developed in this text can be applied to essentially any technology.
References


Reading References


Chapter 3: Data and Control Flow in Systematic Structures

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Sections:
Notation - - - Two Phase Clocks - - - The Shift Register - - - Relating Different Levels of Abstraction - - - Implementing Dynamic Registers - - - Designing a Subsystem - - - Register to Register Transfer - - - Combinational Logic - - - The Programmable Logic Array - - - Finite State Machines - - - Towards a Structured Design Methodology

The process of designing a large-scale integrated system is sufficiently complex that only by adopting some type of regular, structured design methodology can one have hope that the resulting system will function correctly, and not require a large number of redesign iterations. However, the methodology used should allow the designer to take full advantage of the architectural possibilities offered by the underlying technology.

In this chapter we present a number of examples of data and control flow in regularized structures, and discuss the way in which these structures can be assembled into larger groups to form subsystems, and then these subsystems assembled to form the overall system. The design methodology suggested in this chapter is but one of many ways in which integrated system design may be structured. The particular circuit form presented does tend to produce systems of very simple and regular interconnection topology, and thus tends to minimize the areas required to implement system functions. Arrays of pass transistor logic in register to register transfer paths are used wherever possible to implement system functions. This approach tends to minimize power dissipated per unit area, and, with level restoration at appropriate intervals, tends to minimize the time delay per function. The methodology developed is applied in later chapters to the architecture and design of a data processing path and its controller, which together form a microprogrammed digital computer.

Computer architects, who usually design systems in a rather structured way using commercially available MSI and LSI circuit modules, are often surprised to discover how unstructured is the design within those modules. In principle one can use the basic NAND and NOR logic gates described in Chapter 1 to implement combinational logic, build latches from these gates to implement data storage registers, and then proceed to design integrated systems using traditional logic design methodology as applied to discrete devices. Integrated systems are often designed this way at the present time. However, it is unlikely that such unstructured approaches to system design can survive as the technology scales down towards maximum density VLSI.
There are historical reasons for the extensive use of random logic within integrated systems. The first microprocessors produced by the semiconductor industry were fairly direct mappings of early generation central processor architectures into LSI. A block diagram of the Intel 4004, the earliest microprocessor to see widespread commercial application, is illustrated in figure 1a. The actual chip layout of the 4004 shown in Figure 1b indicates the complexity of the LSI implementation of this simple central processing unit. Such LSI systems, directly mapping data paths and control functions appropriate in earlier component technologies, of necessity contained a great deal of random logic. However, the extensive use of random logic results in chip designs of very great geometrical and topological complexity, relative to their logical processing power.

To deal with such complexity, system design groups have often stratified the design problem into architecture, logic design, circuit design, and finally circuit layout, with specialists performing each of these levels of the design. Such stratification often precludes important simplifications in the realization of system functions.

Switching theory provides formal methods for minimizing the number of gates required to implement logic functions. Unfortunately, such methods are of little value in VLSI systems, since the area occupied on the silicon surface by circuitry is far more a function of the topological properties of the circuit interconnections than it is of the number of logic gates implemented. The minimum gate implementation of a function often requires much more surface area for its layout than does an alternative design using more transistors but having simpler interconnection topology.

There are known ways of structuring integrated circuit designs implemented using traditional logic design methods. A notable example is the polycell technique. In this technique, a group of standard cells corresponding to typical SSI or MSI functions are gathered into a library of functions. The logic diagram for the system to be implemented is used to specify which cells in the library are required. The cells are then placed into a chip layout, and interconnections laid out between them by an automatic interconnection routing system. The polycell technique provides the logic designer having limited knowledge of integrated systems with a means of implementing modest integrated circuit designs directly from logic equations. However, a heavy penalty is paid in area, power, and delay time. Such techniques, while valuable expediens, do not take advantage of the true architectural potential of the technology, and do not provide insight into directions for further progress.
The Intel 4004 Microprocessor: An Early LSI System

Fig. 1a. 4004 Block Diagram

Fig. 1b. 4004 Chip Photomicrograph with Pin Designations
Switching theory not only yields the minimum number of gates to implement a logic function, but it also directly synthesizes the logic circuit design. Unfortunately, at the present time there is no general theory which provides us with a lower bound on area, power, and delay time for the implementation of logic functions in integrated systems. Theoretical lower bounds for certain special structures and algorithms of interest are given in chapter 9.

In the absence of a formal theory, we can at best develop and illustrate alternative design methodologies which tend to minimize these physical parameters. Proposed design methodologies should in addition provide means of structuring system designs so as to constrain complexity as circuit density increases. We hope that the examples and techniques presented in this text will serve to clarify these issues and stimulate others to join in the search for more definitive results.

Notation

There are a number of different levels of symbolic representation for MOS circuits and subsystems used in this text. Figures 2a., 2b., 2c., and 2d., illustrate a NAND gate at several such levels. At times it may be necessary to show all the details of a circuit's layout geometry in order to make some particular point. For example, a clever variation in some detail of a circuit's layout geometry may lead to a significant compaction of the circuit's area without violating the design rules.

Often, however, a diagram of just the topology of the circuit conveys almost as much information as a detailed layout. Such stick diagrams may be annotated with important circuit parameters if needed, such as the L/W ratios shown in figure 2b. Many of the important architectural parameters of circuits and subsystems are a reflection of their interconnection topologies.

Alternative topologies often lead to very different layout areas after compaction. The discovery of a clever starting topology for a design usually provides far better results than does the application of brute force to the compression of final layout geometries. For this reason, many of the important structural concepts in this chapter and throughout the text will be represented for clarity by use of colored stick diagrams. The color coding of the stick diagrams is the same as that of layout geometries, and is as follows: green symbolizes diffusion and transistor channel region; yellow symbolizes ion implantation for depletion mode transistors; red symbolizes polysilicon; blue symbolizes metal; black symbolizes a contact.
Later, through a number of examples in chapter 4, we will present the details of procedures by which the stick diagrams are transformed into circuit layouts, and then digitized for maskmaking. Note that if this topological form of representation were formalized, one might consider "compiling" such descriptions by implementing algorithms which "flesh out and compress" the stick diagrams into the final layout geometries\textsuperscript{3}, according to the constraints imposed by the design rules.

When the details of neither geometry nor topology are needed in the representation, we may revert to the familiar circuit diagrams and logic symbols. At times we may find it convenient to mix several levels in one diagram, as shown in figure 2e. A commonly used mixture is: (i) stick diagrams in portions where topological properties are to be illustrated, (ii) circuit symbols for pullups, and (iii) logic symbols, or defined higher level symbols, for the remaining portions of the circuit or system.

We will define logic variables in such a way that a high voltage on a signal path representing that variable corresponds to that variable being true (logic-1). Conversely, a low voltage on a signal path representing that logic variable corresponds to the variable being false (logic-0). Here high voltage and low voltage mean well above and well below the logic threshold of any logic gates into which the signal is an input. This convention simplifies certain discussions of logic variables and the voltages on the signal paths representing them. Thus when we refer to the logic variable $\beta$ being high, we indicate simultaneously that $\beta$ is true (logic-1) and is represented on the signal path named $\beta$ by a high voltage, one well above the logic threshold. In boolean equations and logic truth tables we use the common notation of 1 and 0 to represent true and false respectively, and by implication high and low voltages on corresponding signal paths.
Fig. 2a. NAND Gate: Layout Geometry

Fig. 2b. NAND Gate: Topology (Stick Diagram)

Fig. 2c. NAND Gate: Circuit Diagram

Fig. 2d. NAND Gate: Logic Symbol

Fig. 2e. Example of Mixed Notation
Two Phase Clocks

We will often make use of a particular form of "clocking" scheme to control the movement of data through MOS circuit and subsystem structures. By clocking scheme we mean a strategy for defining the times during which data is allowed to move into and through successive processing stages in a system, and for defining the intervening times during which the stages are isolated from one another.

Many alternative clocking schemes are possible, and a variety are in current use in different integrated systems. The clocking scheme used in an integrated system is closely coupled with the basic circuit and subsystem structuring, and has major architectural implications. For clarity and simplicity we have selected one clocking scheme, namely two-phase, non-overlapping clock signals. This scheme is used consistently throughout the text, and is well matched to the type of basic structures possible in MOS technology.

The two clock signals $q_1$ and $q_2$ are plotted as a function of time in figure 3. The signals both switch between zero volts (logic-0) and a voltage near VDD (logic-1), and both have the same period, $T$. Note that both signals are non-symmetric, and have non-overlapping high times. The high times are somewhat shorter than the low times. Thus $q_2$ is low all during each of those time intervals from when $q_1$ rises, nears VDD, and then falls back to zero.

We have adopted a compact convention for transistions of clocking signals. The rising transition of a signal $q$ is symbolized as \( r_q \), and the fall as \( f_q \). we also have a similar rule for $q_1$, namely $q_1 = 0$ all during each time interval from $r_{q_2}$ to $f_{q_2}$. Therefore, at all times the logic AND of the two signals equals zero: \( [q_1(t)] \cdot [q_2(t)] = 0 \), for all $t$. For convenience, we will often use the following equivalence in our descriptions: "during $q_1$" is equivalent to "during the time period when $q_1$ is high". In the next section we will illustrate the use of these two clocking signals to move data through some simple MOS circuit structures. A more detailed discussion of clocking requirements is given in chapter 7.
The Shift Register

Perhaps the most basic structure for movement of a sequence of data bits is the serial shift register, shown in circuit diagram form in figure 4a. The shift register is composed of level restoring inverters coupled by pass transistors, with the movement of data controlled by applying clock signals \( q_1 \) and \( q_2 \) to the gates of alternate pass transistors in the sequence.

Data is shifted from left to right as follows. Suppose a logic signal \( X \) is present on the leftmost input to the shift register when clock signal \( q_1 \) rises. Then, during the time when \( q_1 \) is high, this signal will propagate through the pass transistor and be stored as charge on the input capacitance of the first inverter stage. For example, if the signal \( X \) is low, then the inverter input gate capacitance will be discharged towards zero volts during the time when \( q_1 \) is high. On the other hand, if \( X \) is high, the inverter input capacitance will charge up towards \( V_{DD} - V_{th} \) during \( q_1 \).

When the clock signal \( q_1 \) falls, the pass transistor becomes an open circuit, isolating the charge on the input of the inverter. The second clock phase is now initiated by the rise of \( q_2 \). During the time interval when \( q_2 \) is high the logic signal \( X \), now inverted, will flow through the second pass transistor onto the gate of the second inverter. This pattern can be repeated an arbitrary number of times to produce a shift register of any length.

Note that since the clock signals do not overlap, the successive pairs of stages of the shift register are effectively isolated from one another during the transfer of data between inverter pairs. For example, when \( q_1 \) is low, and \( q_2 \) is high, all adjacent inverters connected by the \( q_2 \) controlled pass transistors are in the process of transferring data from the left to the right members of the pairs. All these pairs of inverters are isolated from each other by the intervening \( q_1 \) controlled pass transistors which are all open circuits when \( q_1 \) is low.

It is also important to note that the shortest period, \( T \), we can use for clocks controlling such data transfers is determined by the time required to adequately charge or discharge the inverter input gate capacitance through the pass transistor and the preceding stage pullup or pulldown. To this time must then be added an increment of time sufficient to insure that the clocks do not overlap. For more complex systems, the minimum clock period may be estimated as a function of basic circuit parameters as discussed in Chapter 1.

Figures 4b and 4c illustrate the serial shift register using mixed notations. In figure 4b, each inverter circuit diagram has been replaced by its logic symbol. In figure 4c, the pass transistor
Fig. 3. Two Phase Non-Overlapping Clock Signals

Fig. 4a. Shift Register: Circuit Diagram

Fig. 4b. Shift Register: In Mixed Notation
Fig.4c. Shift Register: More Mixed Notation

Fig.5a. Array of Shift Registers

Fig.5b. Shift-Up Register Array
circuit symbols have been replaced by their stick diagrams. When visualizing the inverter, as represented by its logic symbol, in a circuit structure containing mainly stick diagrams, two points should be kept in mind:

(i) The input to the inverter leads directly to the gate, and thus the gate capacitance, of the inverter's pulldown transistor. This input may be used to store a data bit by isolating the charge representing the bit with a pass transistor. Note that the input path will end up on the poly level within the inverter. A contact cut may thus be required to connect the poly gate and the metal or diffusion path on which the signal enters the inverter.

(ii) Since the connection between the source and gate of the inverter pullup transistor requires a connection of all three conducting levels, the inverter output signal may easily be routed out on any one of the three levels.

Identical serial shift registers can be stacked next to each other and used to move a sequence of data words, as shown in figure 5a. The simple structure in figure 5a anticipates the elegant topological simplicity of many important MOS integrated system functions. By connecting the successive inverter stages with diffusion paths, the pass transistors controlled by the clock signals are formed by simply running vertical clock lines in poly. The structure in figure 5a also anticipates another important point: topological simplification often results when control signals flow on lines that are at right angles to the direction of data flow. In this way as many bits as necessary can be processed in parallel with the same control signals.

The example in figure 5a is so rudimentary it is perhaps difficult to visualize the two clock signals as actually containing control information. Let us consider a slightly more complex example, the shift-up register array shown in figure 5b. In this structure, each data bit moving from left to right during $q_2$ has two alternative pass transistor paths through which it can proceed to the next stage: a straight through path, and a path which shifts it up to the next higher row. If the control signal SH is low, then $[q_2 \cdot SH]$ is high, and the straight through pass transistor paths are used during $q_2$. At the same time, $[q_2 \cdot SH]$ is low, thus preventing data flow through the shift-up pass transistor paths. On the other hand, if SH is high, the straight through pass transistors are off and the shift-up pass transistor paths are used during $q_2$, resulting in the entire data word being shifted vertically as well as horizontally. Here the vertical control lines are run in metal, and the pass transistors are selectively formed by crossing the appropriate diffusion paths with short poly lines.
Relating Different Levels of Abstraction

In the discussions in this chapter, we will not have to make extensive calculations of the detailed electrical behavior of the devices and circuits involved in order to analyze the general behavior of digital logic constructed with these devices and circuits. Most of the examples presented in this chapter, and throughout the text, build upon the use of pass transistors coupling inverting logic stages as a means of structuring designs. The general results of chapter one provide the solutions to most device and circuit problems encountered, such as ratio and delay calculations, etc. In most cases, design concepts can be worked out using stick diagrams, and only at the stage of transforming the circuit topology into the detailed circuit layout geometry will these calculations need to be worked out, either by hand or with circuit simulation programs.

It is important to simplify our mental model of integrated circuitry, so as to more quickly and easily analyze or explain the function of a given circuit, and more easily visualize and invent new circuit structures without drifting too far away from physically realizable and workable solutions. Of course, it is a dangerous practice to oversimplify our abstractions of electronic circuit behavior, and there are some nMOS circuits of deceptively simple appearance which have exceedingly complex behavior. However, throughout large portions of digital integrated systems, if the circuit and subsystem design is structured as suggested in this text, an extremely simple mental model of device and circuit behavior will prove adequate to predict circuit and subsystem behavior.

Figure 6a illustrates a simple way of visualizing the operation of successive inverting logic stages coupled by pass transistors. Assume for the moment that any pass transistors in the paths between stages are on. To visualize the time behavior of an inverter, and the effect of the pullup L/W to pulldown L/W ratio, imagine the flow of current from VDD to GND as the flow of a fluid, and the inverter's two transistors as valves. The basis for thinking of the transistors in this manner is the fluid model of their internal behavior, as given in chapter 1. Whether a transistor is on or off depends upon the voltage, and thus upon the charge, on its control gate, and also on its threshold voltage. The upper "valve" is always open, since the pullup transistor is always on. However, the "valve" corresponding to the pulldown transistor may be either open or closed, depending on the amount of charge on its gate.

In figure 6a, the input to inverter-A is a logic-0, so the pulldown of inverter-A is off, and the lower valve is closed. Current is thus diverted to the large charge storage site corresponding to the gate of the pulldown of inverter-B. At this level of diagram we have reverted to the common
Fig. 6a. A Way of Visualizing the Operation of Successive Inverter Stages

Fig. 6b. Successive Inverter Stages: Circuit Diagram and Logic Diagram

Fig. 6c. Successive Inverter Stages Connected Through a Pass Transistor

[ Illustrating effect of the pass transistor "switch"]
convention of positive charge flow from VDD to ground, rather than electron flow from ground to VDD. If sufficient positive charge has flowed onto this gate, corresponding to a high level of fluid in the tank representing the gate capacitance, then the pulldown of inverter-B is turned on, and thus the lower valve of inverter-B is open. If the lower valve in inverter-B is much larger than the upper one, corresponding to a practical pullup to pulldown size ratio, then the pulldown of inverter-B can sink all the source current provided by the pullup. Also, if given sufficient time and if the connecting pass transistor is on, the pulldown can drain off any charge stored on the succeeding inverter's input gate. Thus we can visualize the sequence of inversions of a logic signal propagating through successive inverter stages as an alternation between high and low levels of fluid in the storage tanks. We can also visualize some of the time behavior of the signal propagation: the larger the gate capacitance, the longer it takes to build up enough charge to open the next stage, and the longer it takes to drain charge off the next stage to turn it off.

Figure 6b represents the same physical circuit modelled in figure 6a, but on successively higher levels of abstraction. When analyzing circuit or logic diagrams showing successive inverting logic stages, as in figure 6b, one should keep the model of figure 6a in mind. Whether one is a novice or an expert in integrated system design, it is very helpful to compress the details of any given lower level of abstraction, so as to reduce the complexity of the problems presented at the next higher level, and enable the mind to span problems of larger scope.

We are now able to visualize a very simple model for the pass transistor: it is in fact like a valve, or "switch" in the path between an inverter and the next charge storage site, i.e. the input gate of the next inverter. Figure 6c shows two inverters coupled by a pass transistor, with the pass transistor informally symbolized as a "switch". In the upper diagram of figure 6c, the pass transistor input is a logic-1, and so the "switch" is in the on position, resulting in the output Z being equal to the input X, after a suitable delay time $\Delta t$. Thus during the time the pass transistor gate input $P = 1$, the output $Z(t) = X(t - \Delta t)$. Here $\Delta t$ is some multiple of the transit time, $\tau$, of the inverter pulldown transistor, as discussed in Chapter 1.

In the lower diagram of figure 6c, the pass transistor "switch" is moved to the off position since $P$ is a logic-0. Therefore, according to our model, the valve in the path between the inverters is shut, and the charge, or lack of charge, is isolated in the storage site. Thus, once the pass transistor "valve" is shut, $Z$ remains at a constant value, independent of changes in $X$. In other words, if $P \to 0$, at $t = t_0$, then $Z(t) = Z(t_0)$, for $t > t_0$.\[\]
These simple visualizations of the inverter and the pass transistor will carry us fairly far into LSI subsystem design. Several logic circuits in this chapter are drawn first in stick diagram form, and then informally sketched with pass transistors replaced with "switches", both to clarify the behavior of the circuits involved, and to further demonstrate the applicability of the model.

Implementing Dynamic Registers

Registers for the storage of data play a key role in digital system design. It is interesting to note that a group of adjacent inverters, with their gates isolatable by pass transistors, can be considered a form of temporary storage register. This arrangement is illustrated in figure 7, which shows two levels of symbolism for this dynamic register. Such a register is very simple in structure. It consists of only three transistors per bit position: the pass transistor and the two transistors of the inverter. However, this dynamic form of register will preserve data only as long as charge can be retained on the inverter input gates. Typically dynamic registers are used in situations where the input gate updating control signals are applied frequently. They are ideal in a clocked system in which they are reloaded every clock cycle, as in the shift register.

Suppose we wish to construct a simple register which can be loaded during the appropriate clock phase under the control of a load signal, and which will retain its information through an indefinite number of successive clock periods until it is reloaded using the load signal. A one bit cell for such a register may be constructed using cross coupled inverters in the configuration shown in figure 8. This register cell is still dynamic in form, since it uses charge storage on the gate of the first inverter to preserve its state. However, it need not be loaded on every successive \( \varphi_1 \) as was the simple register in figure 7. The pass transistor leading to it from the preceding stage is switched on only when both \( \varphi_1 \) and LD are high. On any following \( \varphi_1 \) when LD is low, the cell updates itself by the feedback path through the second pass transistor. Figure 9 illustrates a selectively loadable register composed of such cells. One important feature of this type of register is that it provides as output both the true and complemented forms of the stored data. This feature is often useful when the data is to be processed by a following network of combinational logic.

While there are more elaborate forms of dynamic and static registers, the above two forms are sufficient for many of the required data storage applications within integrated systems.
Fig. 7. A Dynamic Register

Fig. 8. A Selectively Loadable Dynamic Register Cell

Fig. 9. A Selectively Loadable Dynamic Register
Designing a Subsystem

The ideas used to construct simple dynamic registers in the preceding section may be applied to the construction of more sophisticated and interesting subsystems. In this section we will describe the design of a stack. The methodology we use for this specific example we will find appropriate for a wide variety of functional subsystems. We first invent a "cell" which implements the most primitive function of the subsystem. This cell dictates a set of "timing" criteria necessary for its proper operation. The cell geometry together with the timing requirements dictates the design of control "circuits" which will surround an array of the basic cells. Once these control circuits are attached to the cell array, and the necessary "interconnections" are made, the entire assemblage constitutes a functional "module" with a well defined "interface" to the next higher level of design. This interface consists of a functional specification, a geometrical specification, and a set of timing requirements for the control inputs, data inputs, and data outputs.

The stack subsystem is commonly called a last-in, first-out (LIFO) stack. It is also known as a pushdown stack, although we will diagram it horizontally rather than vertically. It is a shift register array with three basic operations: during each full clock period (1) we can push in a new data word at one end of the array, pushing all previously entered words one word position further into the array, or (2) we can leave all words in their current position, or (3) we can pop out a word from the end of the array, pulling all previously entered words back out by one word position.

Figure 10a shows the structure of one horizontal row of the stack. Here we have implemented a shift register which can perform the following three operations: shift data left to right, hold data in place, or shift data right to left. There are four control signals used, two of them being active during $q_1$ and two of them being active during $q_2$. The signals $q_1$ and $q_2$ are our familiar two phase, non-overlapping clock signals.

In order for data to be shifted from left to right, the shift right control line (SHR) is driven high during $q_1$, followed by driving the transfer right control line (TRR) high during $q_2$. The bit of data appearing at the left is thus transferred by this operation onto the gate of the first inverter during $q_1$, and thence to the gate of the second inverter during $q_2$. In order for data to be held in place, the signal transfer left (TRL) is driven high during $q_1$ and transfer right (TRR) is driven high during $q_2$, causing the data to recirculate upon itself without shifting. Note that the data can be obtained at any time from the output of the first inverter. However, since new data
may come to the gate of the first inverter during $\varphi_1$, the only safe time to take data out to the left is during $\varphi_2$. The transfer of data from right to left is caused by driving the shift left control (SHL) line high during $\varphi_2$, followed by driving transfer left (TRL) high during $\varphi_1$.

Figure 10b illustrates a possible topological structure of one horizontal row of the stack. There are two horizontal pathways on the diffusion level for shifting bits right or left. The two inverters for one stage of the row are nested between these paths. VDD, GND, and the four control lines run vertically in metal. The four pass transistors required for controlling the movement of data are conveniently implemented by short poly lines which cross the horizontal diffusion tracks at appropriate positions. Note that the entire row is composed of 180° rotations and repetitions of a basic cell containing one inverter.

In a typical implementation of the complete LIFO stack, a number of such rows run parallel to each other in the horizontal direction. The number of rows is equal to the width in bits of the data words involved. The control lines run vertically across the entire stack, perpendicular to the direction of data flow. For data words of any substantial width, the capacitive loading on the control signals would be sufficient to warrant use of super-buffer drivers.

The stack as a whole may be controlled with only two logic signals: one signalling push, and the other signalling pop. The activation of neither of these two signals causes data to recirculate in place, awaiting the next active instruction.

Let us consider how to derive, from push and pop, the control signals for driving the four control lines SHR, TRR, SHL, TRL. A possible scheme is shown in Fig. 10c. We use random logic for this purpose since only a few gates are required to control the large, regular array of circuit cells in the stack. The operation which determines what the stack will do during the subsequent clock phase is brought in on the path labeled OP. It is important to note in the following that only one signal path (OP) is required to bring in both push and pop logic signals, since these are active on mutually exclusive clock phases.

The control scheme is summarized in the timing diagrams in figure 10e. Here we see that holding OP high during $\varphi_2$, followed by low during $\varphi_1$, implements push. Holding OP low during both $\varphi_1$ and $\varphi_2$ causes the data to recirculate in place. Holding OP high during $\varphi_1$, followed by low during $\varphi_2$, implements pop. Thus, the single signal path, OP, is sufficient to carry both stack control signals into the stack.
Fig. 10a. One Horizontal Row of the Stack

Fig. 10b. Topology of One Horizontal Stack Row
Fig. 10c. Generating the Stack Control Signals

Fig. 10d. Stack Geometry and Interconnect Topology
Fig. 10e. Stack control Signal Timing Diagrams

(PULL:
OP: high in phase 2, and
then low in phase 1;
Causes: SHR, not(TRL)

POP:
OP: high in phase 1, and
then low in phase 2;
Causes: SHL, not(TRL))
During $\varphi_1$, the OP signal is fed through the upper pass transistor into the inputs of the two NOR gates $g_1$ and $g_2$. The outputs for these NOR gates are low during this period, since $\varphi'_2$ is high.

If the incoming OP signal is high while $\varphi_1$ is high, then the lower input of NOR gate $g_2$ will be low. Thus when $\varphi'_2$ falls low, the output of $g_2$ will go high, thereby driving SHL high. If the OP signal is instead kept low while $\varphi_1$ is high, then the output of the NOR gate $g_1$ will go high on the fall of $\varphi'_2$, thereby driving TRR high during $\varphi_2$.

During the period when $\varphi_2$ is high and either the shift left (SHL) or the transfer right (TRR) operation is being executed, the signal on the OP line is being stored on the corresponding input gates of the lower two NOR gates, $g_3$ and $g_4$. Thus, if OP is high while $\varphi_2$ is high, a logic-0 is stored on the input of the NOR gate $g_4$, and during the subsequent $\varphi_1$ high period, SHR will be driven high. Conversely, if OP is low while $\varphi_2$ is high, TRL will be driven high during the following $\varphi_1$ high period.

This kind of control scheme recognizes that there must be a null period between any operation and its next occurrence. Control information is taken in during this period and set up for the subsequent operation. The scheme takes advantage of these null periods, when possible, to perform other operations which can be done without conflict. It is an example of a fundamental design technique which can be extended to larger system structures.

When planning the overall architecture of a larger system, it is often useful to represent subsystems, such as the stack, using a higher level of symbolism. To be truly useful, such representations should, in addition to a functional definition, include the topological factors associated with the interconnection points of the subsystem and the geometrical factors of its shape and relative physical dimensions.

A system level sketch of one particular implementation of the stack is shown in figure 10d. Identical driver circuitry is placed along the top and bottom edges of the shift register array. The transfer right and shift left drivers which are set up during $\varphi_1$ (and active during $\varphi_2$) are placed along the top of the shift register array. The transfer left and shift right drivers which are set up during $\varphi_2$ (and active during $\varphi_1$) are placed along the bottom of the array. The OP bit and the clock signals are required on both the top and the bottom of the shift register array.
The integration of this subsystem into a larger integrated system design will require that the data in and out paths be matched to those of subsystems to which the array is connected, and that the $\varphi_1$, $\varphi_2$, and OP signals be available at either the left or right side of the array. By using system level representations that reflect as closely as possible the dimensions and locations of critical signals in all major subsystems, the interactions between topologies and dimensions of the subsystems can be assessed. The feasibility of an overall system architecture can thus be ensured prior to detailed design and layout.
Register to Register Transfer

From an implementation point of view it is often desirable to combine logic steering functions with the clocking of data into registers, since both require pass transistors as their elementary functional unit. An example is the shift-up register array shown in figure 6. From the next higher level system view, however, it is desirable to separate the two functions conceptually. In Fig. 11a we have shown some combination of inputs, X0 through Xn going through some combination of pass transistors, which may or may not have logic functions attached, into the input gates of some inverting logic elements. This combination of function is then abstracted into a register clocked on the phase during which the input pass transistors are turned on. Any logic function associated with the input pass transistors is considered part of the preceding combinational logic module. This viewpoint is an extension of the concept of dynamic register previously developed in figure 7.

Using this notation, any processing function can be built up using blocks of the form shown in Fig. 11b. Here we have a clocked input register, a block of strictly combinational logic with no timing attached, and an output register clocked on the opposite phase. In this case the inputs are stored in the input register during $q_1$. They then propagate into and through the combinational logic (C/L), with the resulting outputs stored in the output register during $q_2$. Any single data processing step can be viewed as a transfer from one such register to a second through a combinational logic block.

A sequence of such operations can be performed on a data stream by a series of such combinational blocks separated by registers as shown in Fig. 11c. Since different sets of data words in the stream may be operated upon at the same time, but at different locations, this data path is a type of pipelined processing structure. Such pipelined processing structures offer the opportunity for improved processing bandwidth by performing many different operations concurrently. Notice that the throughput rate of such a pipeline system of register to register transfer operations is limited by the delay time through the slowest of the combinational logic blocks. If no registers had been interposed between the function blocks, and each operand set separately run through the entire sequence of combinational logic modules, the throughput rate would be much lower.

In line with the ideas developed earlier in this chapter, the detailed functions performed by the combinational logic modules may often be implemented in circuit structures of very simple and
regular topology. Control signals will in general cross the data path at right angles to the direction of data flow. Figure 11c illustrates sets of such control inputs as \( n_1 \) lines carrying the control function \( OP_1 \) into the first C/L module, \( n_2 \) lines carrying \( OP_2 \) into the second, etc.

The idea of data being processed while passing through combinational logic interspersed between register stages in a sequence of register to register transfers is a basic and important concept in the hierarchy of digital system architecture. We have already described the implementation of registers. The next sections will describe some ways to implement combinational logic functions.

**Combinational Logic**

Combinational logic modules contain no data storage elements. The outputs of a combinational logic module are functions only of the inputs to that module, provided that sufficient time has been allowed for those inputs to propagate through the module's circuitry.

In integrated systems, combinational logic design problems will typically fall within one of three general classes. The first is when a small amount of simple logic is required, for example to derive control signals at the periphery of a system module (as in the stack control signal generation) or to implement a simple function within a single circuit cell (which may then be replicated in a regular array). In these cases, traditional logic design procedures using static NAND and NOR gates can be applied. Such designs involving a few gates are usually rather simple, and can be produced by inspection rather than by use of formal minimization and synthesis procedures. Even in these simple cases, the minimum static logic gate implementation does not necessarily result in either the most regular, the minimum area, the minimum delay, or the minimum power design. In fact, we often find alternative techniques to the use of static logic gates, which in specific instances lead to "better" designs by one of these measures than would minimum gate implementations. For example, figure 12a shows a selector logic circuit (I. Sutherland), in which one of the inputs \( S_1, S_2, S_3, S_4 \) is selected for output by the control variables \( A, B \) according to the function:

\[
Z = S_1A'B' + S_2AB' + S_3A'B + S_4AB
\]

This selector circuit is composed simply of poly paths crossing diffusion paths. Where depletion
Fig. 11a. A Register

Fig. 11b. A Section of a Data Path

Fig. 11c. General Form for a Data Path
mode transistors are placed, the diffusion level path is always connected, thus placing control in the selectively located enhancement mode pass transistors, which function as simple switches. Figure 12c shows the circuit's paths from inputs to outputs using the "switch" abstraction for each of the pass transistors. For each possible combination of values of A and B, there is a path through the selector to Z from only one of the inputs $S_i$. For the specific inputs shown in the example in figure 12c, the signal $S_2$ propagates through to Z since both A and B' are high. Note that no static power is consumed by the circuit, and the area occupied by the circuit is minimal since no contact cuts are required within it. In chapter 5 we describe a very general and powerful arithmetic logic unit (ALU) which uses an array of such selector blocks to control a pass transistor carry network.

The second general class of combinational logic design problems are those rather complex functions for which clever ways of structuring topologically regular implementations have been discovered. As an example, consider the implementation of a tally function. This function has n inputs and n+1 outputs. The kth output is to be high, and all other outputs low, if k of the inputs are high. The boolean equations representing this function for the simple case of three inputs are:

\[
Z_0 = X_1'X_2'X_3' \\
Z_1 = X_1X_2'X_3' + X_1'X_2X_3' + X_1'X_2'X_3 \\
Z_2 = X_1X_2'X_3' + X_1X_2'X_3 + X_1'X_2X_3 \\
Z_3 = X_1X_2X_3
\]

If this function were designed with random logic consisting of active pullup, static logic gates, it would result in a topological kludge. Figure 12b shows a topologically regular implementation of the tally function. A major portion of the function is implemented using a regular array of identical cells each containing only two pass transistors. The design is based on the shift-up register idea presented earlier. A high signal propagates through the array from the pullup at the lower left. Whenever one of the variables $X_i$ is high, the propagating high signal moves up to the next higher horizontal diffusion level path. Thus the number of paths it moves up equals the number of inputs $X_i$ which are high. Logic-0 signals propagate through the array from the ground points to all other outputs.
Figure 12d shows the paths from inputs to outputs for this tally circuit, using the "switch" abstraction for the pass transistors. The figure shows a specific example of a set of inputs controlling the pass transistors of the circuit. Since two of the inputs are high, the logic-1 signal is shifted up two rows and emerges at $Z_2$.

This tally function design can be easily expanded to handle more than three inputs by simply extending the array structure upwards and to the right. However, remember that the delay through $n$ pass transistors is proportional to $n^2$. Thus it may be necessary to insert level restoration prior to such extension. Similar comments apply to the extension of the selector circuit previously shown, or other pass transistor logic arrays one might invent.

The electronic logic gates traditionally used in digital design are unilateral elements: they allow a logic signal to propagate in one direction only. It should be noted that the pass transistor is a bilateral circuit element. It permits the flow of current, and thus the passage of a logic signal, in either direction when its gate is high. While this property of the pass transistor is not necessarily of fundamental importance in integrated systems, it is an interesting and occasionally useful one.

Early relay switching logic used switching contacts which were bilateral elements. Interesting discussions of relay switching logic are contained in both references R4 and R5. The tally array example just given is a basic symmetric network mapped directly into nMOS from relay switching logic (see R5, p.241). The mathematics of switching universally used in digital systems today was proposed by Claude Shannon (R7) in 1938. Shannon demonstrated that the calculus of propositions, based on the algebra of logic developed by Boole (R8), was directly applicable to relay switching circuits.

A third combinational logic design situation occurs when a complex function must be implemented for which no direct mapping into a regular structure is known. Methods for handling this situation are the subject of the next section.

In the design methodology developed in this text, the combinational logic between stages in the register to register transfer paths is often done by operations on the charge moving between stages, using pass transistors to perform these operations. Many researchers at the present time are searching for alternative structures and techniques for performing elementary logic functions, including the use of charge transfer devices$^5$. 


Fig. 12a. Selector Logic Circuit

Fig. 12b. A Tally Circuit
Fig. 12c. Example of Operation of Selector Circuit

Fig. 12d. Example of Operation of Tally Circuit
(visualizing where the switches are)
The Programmable Logic Array

On many occasions it is convenient to implement the combinational logic interspersed between register stages with regular structures of pass transistors. However, we will often encounter important combinational logic functions which do not map well into such regular structures. In particular, combinational logic used in the feedback paths of finite state machines is often highly complex and inherently irregular. Also, we may wish to delay binding the details of the logic functions used in finite state machine sequencing until most of the design is complete. If the combinational logic were implemented in an irregular structure, such changes could require a major redesign.

Fortunately, there is a way to map irregular combinational functions onto regular structures, using programmable logic arrays (PLA's) as described in this section. This technique of implementing combinational functions has a great advantage: functions may be significantly changed without requiring any major design or layout changes of the PLA structure.

One very general and regular way to implement a combinatorial logic function of $n$-inputs and $m$-outputs is to use a memory of $2^n$ words of $m$-bits each. The $n$-inputs form an address into the memory, and the $m$-outputs are the data contained in that address. Such a memory implements the full truth table for the output functions. Many systems are in fact built using memories as combinational logic elements. A common form of memory for this purpose is the read-only memory (ROM) where the data is permanently placed in the memory by a mask pattern, or by electrically altering the individual bit positions. There is one major difficulty with this approach: it is often the case that most of the possible input combinations cannot occur, due to the nature of the specific problem. Stated another way, many combinational logic functions require only a small fraction of all $2^n$ product minterms for a canonical sum of products implementation. In such cases, a ROM is very wasteful of area.

The programmable logic array (PLA) is a structure which has all the generality of a memory for implementing combinational logic functions. However, any specific PLA structure need contain a row of circuit elements only for each of those product terms that are actually required to implement a given logic function (see R4, Ch.4). Since it does not contain entries for all possible minterms, it is usually far more compact than a ROM implementation of the same function. To achieve full compaction, the various output functions must be jointly minimized before the PLA layout pattern can be defined. However, such minimization is not essential. Less than full
compaction increases the independence of the different entries, so that changes in function may require only local changes in the PLA.

An illustration of the overall structure of a PLA is shown in figure 13a. The diagram includes the input and output registers, in order to show how easily these are integrated into the PLA design. The inputs, stored during \( q_1 \) in the input register, are run vertically through a matrix of circuit elements called the AND-plane. The AND-plane generates specific logic combinations of the inputs and their complements. The outputs of the AND-plane leave at right angles to its inputs and run horizontally through another matrix called the OR-plane. The outputs of the OR-plane then run vertically and are stored in the output register during \( q_2 \).

The circuit diagram of a specific programmable logic array is shown in figure 13b. This diagram will help to clarify the structure and function of the AND and OR-planes of the PLA. The input register bit for each input path is formed by a pass transistor clocked on \( q_1 \) leading to both inverting and non-inverting super buffers. These buffers drive two lines running vertically through the AND-plane, one for the input term and one for its complement. The outputs of the AND-plane are formed by horizontal lines with pull-up transistors at their leftmost end. The function of the PLA's AND-plane is then determined by the locations and gate connections of pull-down transistors connecting the horizontal lines to ground.

Each output running horizontally from the AND-plane carries the NOR combination of all input signals which lead to the gates of transistors attached to it. For example, the horizontal row labelled \( R_3 \) has three transistors attached to it in the AND-plane, one controlled by A, one by B and one by C'. If any of these inputs is high, then \( R_3 \) will be pulled down towards ground and will be low.

Thus, \( R_3 = (A + B + C')' = A'B'C. \) Similarly, \( R_4 = (A + B' + C)' = A'BC'. \)

The OR-plane matrix of circuit elements is identical in format to the AND-plane matrix, but rotated 90 degrees. Once again, each of its outputs is the NOR of the signals leading to the gates of all transistors attached to it. In figure 13b for example, both \( R_3 \) and \( R_4 \) lead to the gates of transistors leading from the output line \( Z_4' \) to ground. If either \( R_3 \) or \( R_4 \) is high, \( Z_4' \) will be low. Thus, \( Z_4' = \text{NOR}(R_3, R_4) = (A'B'C + A'BC')' \). Up to this point the PLA implements the \textit{NOR-NOR canonical form} of boolean function of its inputs.

The output lines of the OR-plane matrix are run into an output register formed by pass
Fig. 13a. Overall Structure of the PLA

Fig. 13b. Circuit Diagram of PLA Example
Fig. 13c. Stick Diagram of PLA Example

**Product Terms:**

\[
\begin{align*}
R_1 &= (A')' = A \\
R_2 &= (B+C)' = B'C' \\
R_3 &= (A+B+C)' = A'B'C \\
R_4 &= (A+B'+C)' = A'BC'
\end{align*}
\]

**Outputs:**

\[
\begin{align*}
Z_1 &= A \\
Z_2 &= A + A'B'C \\
Z_3 &= B'C' \\
Z_4 &= A'B'C + A'BC'
\end{align*}
\]
transistors (clocked on $\phi_2$) leading into inverting drivers. Note that the output $Z_4$ at this point is: $Z_4 = A'B'C + A'BC$. This expression illustrates why the two PLA planes, each implementing the NOR function, are usually referred to as the AND and OR-planes. Following the output register, the outputs appear directly as the sum of products canonical form of boolean functions of the PLA inputs, that is, the OR of AND terms. Each horizontal line of the PLA carries one product term.

Figure 13c shows one possible layout topology for implementing the PLA in nMOS circuitry. The example is the same circuit illustrated in figure 13b. The input lines crossing each plane are run in poly. The output lines from each plane are run in metal. Paths running to ground are placed between alternate poly lines, on the diffusion level. It is then a simple matter to form the pulldown transistors connecting the metal output lines to ground. They are selectively located diffusion lines under the appropriate input poly lines.

Although the PLA may implement a very irregular combinational function, the irregularity is confined to the irregular locations of pulldown transistors which "program" the function. The overall structure and topology of the PLA are very regular. Note that its overall shape and size is a function of the parameters: (i) the number of inputs, (ii) the number of product terms, (iii) the number of outputs, and (iv) the length unit $\lambda$.

Finite State Machines

In many cases in the processing of data, it is necessary to know the outcome of the current processing step before proceeding with the next. Results of the current step may be used as inputs in the next step. The configuration shown in figure 14a can be used to implement a processing stage having this requirement. A typical register to register transfer stage has been modified by simply feeding back some of its outputs to some of its inputs. This structure implements a form of sequential machine known as a finite state machine.

The feedback signals form a binary number which may be regarded as identifying the state of the machine. The value of this number is stored, along with the external inputs, in the first register during $\phi_1$. These combined inputs then propagate through the combinational logic. The resulting outputs are stored in the second register during $\phi_2$. The falling edge of $\phi_2$ must occur a sufficient time later to insure that all signals have propagated through the combinational logic.
Each complete machine cycle, consisting of $q_1$ followed by $q_2$, results in two new sets of outputs: (i) the external outputs which are typically used for controlling other units of the system, and (ii) a new feedback number, which defines the next state of the machine. This process repeats during each clock period. The number of possible states is determined by the number of bits in the feedback path, and is finite.

There are a number of ways of abstractly representing the states, the required state transitions, and the outputs of sequential machines under given input sequences. Possible representations include state diagrams, transition tables, boolean or numerical difference equations, etc. A large body of theory has been developed concerning sequential machines. The serious reader will benefit from a further study of the results of switching theory on this subject (R3, R4).

Implementations of simple finite state machines are used to produce the very lowest level of system control sequencing, since they can autonomously generate control sequences. The sequential machine having a finite number of states is a very important element in the hierarchy of fundamental concepts used in integrated system architecture.

The configuration shown in figure 14a implements a synchronous machine, since the feedback loop is only activated at times determined by the clock signals. In any clock period $k$, the output terms $Z_j$ and the next state terms $Y_f$ are valid during $q_1(k)$. They are functions of the external inputs $X_i$ and feedback terms $Y_f$ which were valid during $q_1(k-1)$.

If a sequential machine contains a feedback loop which is continuously active, then it may begin a response to a change in inputs or state at any time, rather than just at fixed clock times. Such a sequential machine is referred to as an asynchronous sequential machine. The analysis of asynchronous sequential machines and their implementation is far more complex than that of synchronous ones. Great care must be exercised to avoid any difference in state sequencing and outputs under arbitrary differential delays of signals through the circuit paths of such machines (R3, Ch.5). There will be only a few special cases where we use the asynchronous form of sequential machine (Chapter 7), and these will be subject to detailed analysis.

Where sequential machines are required within integrated systems, we will generally implement them in synchronous form. Synchronous machines are rather easy to implement correctly, and fit naturally into the two phase clocking scheme used for moving data around within our systems.

However, the reader should carefully note that an implementation of a synchronous sequential
Fig. 14a. Feedback in Register Transfer Path, Implementing a Finite State Machine

Fig. 14b. PLA Implementation of a Finite State Machine
machine functions correctly only if the delays in the circuit paths are sufficiently short compared to the clock period. If we were to implement many copies of a particular machine, the probability of correct function for any given copy is thus a function of (i) the clock period used, and (ii) the distribution of differential delays in that copy's signal paths. Our estimate that a particular copy will function correctly is thus based in part on assumptions about the ratio of likely deviations in circuit delays to the clock period. A discussion of delays in MOS circuits is given in chapter 1.

There is a very straightforward way to implement simple finite state machines in integrated systems: we use the PLA form of combinational logic and just feed back some of the outputs to the inputs, as illustrated in figure 14b. The circuit's structure is topologically regular, has a reasonable topological interface as a subsystem, and is of a shape and size that are functions of the appropriate parameters. The function of this circuit is determined by the "programming" of its PLA logic. If, for example, early in a design cycle there is some uncertainty in the details of the desired sequencing of such a circuit, it is easy to provide layout space for extra, unused inputs, minterms, or outputs as contingencies.

The following simple example will help illustrate the basic concepts of finite state machines and their implementation in nMOS circuitry. A busy highway is intersected by a little used farmroad, as shown in figure 15a. Detectors are installed which cause the signal C to go high in the presence of a car or cars on the farmroad at the positions labelled C. We wish to control traffic lights at the intersection, so that in the absence of any cars waiting to cross or turn left on the highway from the farmroad, the highway lights will remain green. If any cars are detected at either position C, we wish the highway lights to cycle through caution to red, and the farmroad lights then to turn green. The farmroad light is to remain green only while the detectors signal the presence of a car or cars, but never longer than some fraction of a minute. The farmroad light is then to cycle through caution to red, and the highway light then to turn green. The highway light is not to be interruptible again by the farmroad traffic until some fraction of a minute has passed.

A state diagram model of a finite state machine to control the lights is sketched in figure 15b. This diagram identifies four possible states of the machine, and indicates the input conditions which cause all possible state transitions. A block diagram of the PLA circuit implementing the machine is shown in figure 15c. The circuit uses the signal C as an input, and provides outputs HL and FL which encode the colors of the highway and farmroad lights it controls. Note that a
timer is used to provide, as controller inputs, the short and long timeout signals (TS, and TL), at appropriate times following a start timer (ST) signal output from the controller. This timer could be implemented as a digital counter in the same nMOS circuitry. Another abstract model describing the desired function of the controller is given in the state transition table in figure 15d, which contains similar information to that in the state diagram.

The detailed sequencing of the machine under various input sequences is described by both the state diagram and transition table models of the controller. Consider starting in the state HG, where the highway lights are green. The machine remains in state HG as long as either no cars are detected or the long timeout has not occurred, in other words as long as \((C \text{AND} TL) = 0\). After the long timeout occurs, if any cars are detected, the machine restarts the timer and changes state to HY, where the highway lights are yellow. It remains in state FY only until the short timeout occurs, and then restarts the timer and changes to state FG, where the farmroad lights are green. It remains in state FG until either no cars are detected or the long timeout occurs, i.e. \((C \text{OR} TL) = 1\). Then it restarts the timer and changes to state FY, where the farmroad light is yellow. It remains in state HY only until the short timeout occurs. It then restarts the timer and changes to state HG, the starting state.

The locations of transistors in the PLA light controller circuit can be determined by "hand assembling" the "program" specified in the "symbolic" transition table in figure 15d, resulting in the encoded state transition of figure 15e. First we assign codes to the states. In the example: state HG is encoded as \((Y_0, Y_1) = (0,0)\); HY as \((0,1)\); FG as \((1,1)\); and FY as \((1,0)\). Next, we assign codes to the output light control signals: green is encoded as \((0,0)\), yellow as \((0,1)\), and red as \((1,0)\). We now form the encoded state transition table by constructing one row for each product term implied by the symbolic table of fig. 15d. A row in 15d specifying a state transition as a function of a single input variable or single product term of input variables produces a single row in table 15e. A row in table 15d specifying a state transition as a function of a sum or sum of products of input variables, leads to a corresponding number of rows in table 15e.

Placement of the transistors within the PLA matrices follows directly from the encoded state transition table:

(i) For each logic-1 in the next state and output columns in the table, we run a diffusion path from the corresponding next state or output line in the PLA OR-plane, under the corresponding product term line, to ground. This creates a transistor controlled by the product term line. Then,
Fig. 15a. A Highway Intersection

Fig. 15b. Light Controller State Diagram

Fig. 15c. Controller Block Diagram

Fig. 15d. Transition Table for the Light Controller

<table>
<thead>
<tr>
<th>In Present State:</th>
<th>If Inputs are:</th>
<th>Next State will be:</th>
<th>and Outputs are:</th>
</tr>
</thead>
<tbody>
<tr>
<td></td>
<td></td>
<td></td>
<td>HL</td>
</tr>
<tr>
<td>Highway Green</td>
<td>(Cars \land (TimeoutL) = 0</td>
<td>Highway Green</td>
<td>Green</td>
</tr>
<tr>
<td></td>
<td>(Cars \land (TimeoutL) = 1</td>
<td>Highway Yellow</td>
<td>Green</td>
</tr>
<tr>
<td>Highway Yellow</td>
<td>TimeoutS = 0</td>
<td>Highway Yellow</td>
<td>Yellow</td>
</tr>
<tr>
<td></td>
<td>TimeoutS = 1</td>
<td>Farmroad Green</td>
<td>Red</td>
</tr>
<tr>
<td>Farmroad Green</td>
<td>(Cars') \lor (TimeoutL) = 0</td>
<td>Farmroad Green</td>
<td>Red</td>
</tr>
<tr>
<td></td>
<td>(Cars') \lor (TimeoutL) = 1</td>
<td>Farmroad Yellow</td>
<td>Red</td>
</tr>
<tr>
<td>Farmroad Yellow</td>
<td>TimeoutS = 0</td>
<td>Farmroad Yellow</td>
<td>Red</td>
</tr>
<tr>
<td></td>
<td>TimeoutS = 1</td>
<td>Highway Green</td>
<td>Red</td>
</tr>
</tbody>
</table>

* Inputs not listed = don't cares
Fig. 15e. Encoded State Transition Table for the Light Controller

<table>
<thead>
<tr>
<th>Inputs:</th>
<th>Present State:</th>
<th>Next State:</th>
<th>Outputs:</th>
<th>Product terms:</th>
</tr>
</thead>
<tbody>
<tr>
<td>C TL TS</td>
<td>( Y_{p0}, Y_{p1} )</td>
<td>( Y_{n0}, Y_{n1} )</td>
<td>ST HL0 HL1 FL0 FL1</td>
<td>R1 R2 R3 R4 R5 R6 R7 R8 R9 R10</td>
</tr>
<tr>
<td>0 X X</td>
<td>0, 0 (HG)</td>
<td>0, 0 (HG)</td>
<td>0 0 0 1 0</td>
<td></td>
</tr>
<tr>
<td>X 0 X</td>
<td>0, 0 (HG)</td>
<td>0, 0 (HG)</td>
<td>0 0 0 1 0</td>
<td></td>
</tr>
<tr>
<td>1 1 X</td>
<td>0, 0 (HG)</td>
<td>1, 0 (HY)</td>
<td>1 0 0 1 0</td>
<td></td>
</tr>
<tr>
<td>X X 0</td>
<td>0, 1 (HY)</td>
<td>0, 1 (HY)</td>
<td>0 0 1 1 0</td>
<td></td>
</tr>
<tr>
<td>X X 1</td>
<td>0, 1 (HY)</td>
<td>1, 1 (FG)</td>
<td>1 0 1 1 0</td>
<td></td>
</tr>
<tr>
<td>1 0 X</td>
<td>1, 1 (FG)</td>
<td>1, 1 (FG)</td>
<td>0 1 0 0 0</td>
<td></td>
</tr>
<tr>
<td>0 X X</td>
<td>1, 1 (FG)</td>
<td>1, 0 (FY)</td>
<td>1 1 0 0 0</td>
<td></td>
</tr>
<tr>
<td>X 1 X</td>
<td>1, 1 (FG)</td>
<td>1, 0 (FY)</td>
<td>1 1 0 0 0</td>
<td></td>
</tr>
<tr>
<td>X X 0</td>
<td>1, 0 (FY)</td>
<td>1, 0 (FY)</td>
<td>0 1 0 0 1</td>
<td></td>
</tr>
<tr>
<td>X X 1</td>
<td>1, 0 (FY)</td>
<td>0, 0 (HG)</td>
<td>1 1 0 0 1</td>
<td></td>
</tr>
</tbody>
</table>

Fig. 15f. PLA Sequential Circuit Implementing the Light Controller
if that controlling product term line is ever high, the path to the output inverter will be low, and the output will be high. The output line will be low unless some product term line controlling it is high.

(ii) For each logic-1 in the input and present state columns in the table, we run a diffusion path from the corresponding product term line, under the corresponding inverted input or state line in the PLA AND-plane, to ground. The transistor thus created is controlled by the inverted input or state line. Whenever that controlling line crossing the AND-plane is high, the product term line will be low.

(iii) For each logic-0 in the input and present state columns in the table, we run a diffusion path from the corresponding product term line, under the corresponding non-inverted input or state line in the PLA AND-plane, to ground. The transistor thus created is controlled by the non-inverted input or state line. Whenever that controlling line crossing the AND-plane is high, the product term line will be low.

Note that if all lines which control the transistors connecting a given product term line to ground are low, then that product term line will be high. Otherwise it will be low.

The PLA circuit in figure 15f is programmed from the transition table in figure 15e, according to the rules above, and implements the traffic light controller. Note that this LSI implementation does not exactly strain itself to meet the time response requirements of the control problem: it can run at a clock rate at least 10^7 times as fast as required. Also, note that the PLA controller is roughly (150Λ)^2 in area. Using the 1978 value of Λ = 3μm, this controller is (450μm)^2 ~ 0.002 cm^2 in area. A PLA controller this size may contain over 150 transistors, but occupies only 1/125th of the area of a typical 0.25 cm^2 silicon chip in 1978. By the late-80's, as Λ scales down towards its ultimate limits, such a controller will require only ~ 1/25,000th of the area of such a chip.

As we will see in later chapters, a data processing machine of any desired complexity can be created by interconnecting register to register data processing paths constructed along the lines of that shown in figure 11c, such paths being controlled by finite state machines implemented as shown in figure 14b. The data paths form the "highways" for the movement of data, under control of the finite state machine "traffic controllers".
Towards a Structured Design Methodology

The task of designing very complex systems involves managing, in some highly structured way, the space and time relationships between the various levels of system building blocks so that the entire system will function as intended when it is finished. The beginnings of a structured design methodology for VLSI systems can be produced by merging together in a hierarchy the concepts presented in this chapter. Designs are then done in a "top down" manner, but with a full understanding by the architect of the successive lower levels of the hierarchy.

To begin, we plan our digital processing systems as combinations of register to register data transfer paths, controlled by finite state machines. Then the geometric shapes, relative sizes, and interconnection topologies of all subsystem modules are collectively planned so all modules will merge together snugly, with a minimum of space and time wasted by random interconnect wiring. Storage registers are typically constructed by using charge stored on input gates of inverting logic. The combinational logic in the data paths is typically implemented using steering logic composed of regular structures of pass transistors. Most of the combinational logic in the finite state machines is typically implemented using PLA's. All functioning is sequenced using a two-phase, non-overlapping clock scheme.

When viewed in its entirety, a system designed in this manner is seen as a hierarchy of building blocks, from the very lowest level device and circuit constructs, on up to and including the high level system software and application programs in which the intended functions of the system are finally expressed. Individuals who understand the key concepts of each level in this hierarchy will recognize that the boundaries between levels are rather elastic ones. Each level of activity might best be optimized not on its own as a specialty, but as it fits into an overall system's picture. For example, the activity "logic design" in integrated systems might best be conceptualized as the search for techniques and inventions which best couple the physical, topological, and geometric properties of integrated devices and circuits with the desired properties of digital VLSI systems. The search for alternative components for any given design hierarchy, and the search for alternative hierarchies, will be done best by those who span more than one specialty.

A particularly uniform view of such a system of nested modules emerges if we view every module at every level as a finite state machine or data path controlled by a finite state machine. At the lowest level, elements such as the stack and register cells may be viewed as state machines with
one feedback term (the output), two external inputs (the control signals), and a one bit state register. These rudimentary state machines are grouped in a structured manner to form portions of a state machine, or data path controlled by a state machine, at the next level of the hierarchy. Structured arrays of identical state machines often provide a mechanism for distributing processing among memory cells, thus enabling vast increases in processing bandwidth. Although in some cases the feedback paths are used in rather specialized ways, the state machine metaphor still provides a precise description of module behavior. The entire system may thus be viewed as a giant hierarchy of nested machines, each level containing and controlling those below it. A detailed quantitative treatment of certain hierarchically organized machines is given in chapter 9.

In chapters 5 and 6 we will apply the design methodology developed in this chapter to the design of a digital computer system. A one chip implementation of the data path portion of this computer system is illustrated in the frontispiece. Consistent use of the described design methodology resulted in a design of great regularity, short delay times, low power consumption, and high logical processing capability. As we will see in chapter 4, regular designs, with small numbers of basic circuit cell types replicated in two dimensions to form subsystems, also have significant implementation advantages over less structured designs.

References

1. - - - polycell reference: - - -


Reading References


R2. B. Sucek, "Microprocessors and Microcomputers", John Wiley, 1976, is a good introductory reference containing sections on basic digital design, and on the interfacing and programming of a number of present day microprocessors.


Chapter 4: Implementing Integrated System Designs:
From Circuit Topology to Patterning Geometry to Wafer Fabrication
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Sections:

Patterning and Fabrication • Hand Layout and Digitization using a Symbolic Layout Language • An Interactive Layout System • The Caltech Intermediate Form for LSI Layout Description • The Multi-Project Chip • Examples • Patterning and Fabrication in the Future • Fully Integrated, Interactive Design Systems • System Simulation, Test Generation, and Testing

This chapter presents the basic concepts involved in implementing integrated system designs, from the system designer's point of view. Tools are described which help the designer produce the geometrical layout patterns for each layer of an integrated system given the logic, circuit, or topological level design of the system. Procedures are described for encoding these layout patterns and then using the encoded layouts in the patterning and fabrication processes to implement the integrated system. In addition, we discuss how design tools and procedures are likely to evolve towards fully integrated design systems, under the influence of increased complexity of design and predictable changes in the technologies of implementation.

To enable groups of readers to actually design moderate sized LSI systems, we've included descriptions of easily constructed LSI design tools and procedures for organizing and implementing LSI multi-project chips. In each case, the tools are described as part of a complete system of design and implementation procedures, some of which are performed manually while others are machine assisted. Those experienced in software system design will recognize that construction of the machine-assisted portions of these systems is fairly straightforward. Contrary to what many may think, designing your own LSI projects, merging them onto collaborative multi-project chips, and having these implemented by commercial maskmaking and wafer-fabrication firms is now well within the computational and financial reach of most industrial R&D groups and university EE/CS departments.

We are firm believers in learning by doing, and hope that the information provided in this chapter will help and encourage many groups of readers to try their hand at building LSI design tools and designing LSI systems. Such first-hand experience will lead to a deeper understanding of the remaining material in this text.
An overview of the stages of integrated system design, layout, and implementation is given in figure 1. The designer first transforms the circuit and topological level designs into a geometrical layout of the system, using procedures described later in this chapter. In order to optimize the layout, perform various design checks, and discover errors, the designer usually iterates several times between design and layout. The result is a set of design files describing the layout. These files are in a particular representation called an intermediate form, which efficiently and unambiguously describes the layout geometry.

The design files are then converted into files for driving the chosen patterning mechanism. At present, design files are commonly converted into pattern generator (PG) files, for use by a maskmaking firm for driving an optical pattern generator, the first step of maskmaking. By a sequence of photolithographic steps, the mask house produces a set of masks, which a commercial wafer fabrication firm may then use to pattern silicon wafers. Each finished wafer contains an array of system chips. The wafers are then diced into separate chips, which are packaged and tested to yield working systems.

From the system designer's point of view, maskmaking and fabrication can be visualized as one would a film processing service: the designer produces the "artwork" (design files), from which the mask house makes "negatives" (masks), which are then run on a fab line to produce "prints" (wafers). The maskmaking and fabrication sequence is function, design, and layout independent: the mask and fab firms do not require detailed information about the integrated systems they fabricate. If the original layouts satisfy the design rules, and satisfy a few constraints imposed by patterning and fabrication, then these processes will yield correctly patterned wafers.

One need not closely bind a system's design to the detailed processing specifications of particular mask and fab firms. Various firms will differ somewhat in the minimum value of the length unit $\lambda$ which they can successfully process. The transit time of the transistors fabricated, and the resistance per square and capacitance per unit area of fabricated features will also vary from one fab line to another. However, well structured and relatively process independent nMOS designs will function correctly if scaled to a value of $\lambda$ appropriate for the chosen fabrication facilities, and operated using an appropriate system clock period.

We next examine some of the present implementation procedures a bit more closely, to set the stage for sections on design and layout. Those later sections will be clearer if one can visualize how the design files are to be used during patterning and fabrication.
Fig. 1. Overview of Integrated System Implementation
Patterning and Fabrication

On completion of design and layout, the system design is contained in system layout files in intermediate form. Prior to fabrication, a final check plot of the layout is usually generated by converting these design files into files for driving a graphics plotter. Check plots are used for visually checking for design rule violations and other design errors. Once the designers have done as much visual checking as they are going to do, the system layout files are converted into pattern generator (PG) files, to be sent to the maskmaking facility. Figure 2 summarizes the sequence of patterning and fabrication procedures which then follows, and identifies the artifacts passed on at each step in the sequence.

Maskmaking begins with pattern generation to produce reticles. Present pattern generators are projector-like systems containing (i) a precisely movable stage, (ii) an aperture of precisely variable rectangular size and angular orientation, and (iii) a light source, all program controllable by a computer system. To produce a reticle, a photographic plate is mounted on the stage, and the PG file for a particular system layer is used to direct the flashing of a sequence of rectangular exposures, of particular sizes and orientations, onto a sequence of coordinate locations on the plate, as illustrated in figure 3.

The PG file contains a sequence of entries, each of which describes a rectangle\(^5\). A typical representation uses five numbers for each rectangle: the \(x,y\) coordinates of its center, and its height, width, and angular orientation, as shown in figure 4. One can now visualize the nature of the conversion from intermediate form to PG files: the layout of each layer must be decomposed into its equivalent as a set of rectangles, each having \([x,y,h,w,a]\) values flashable by the particular pattern generator, and these rectangles must be sorted into an efficient flashing sequence for that pattern generator.

When the flashing sequence is completed, the plate is developed, yielding the reticle. A sketch of such a reticle is given in figure 5. Each reticle is a photographic master copy much like a photo negative, of the layout of one system layer, usually at a scale ten times (10x) the final system chip size. Photo enlargements of reticles, called "blowbacks", may be obtained from the mask house, to provide a further level of checking of design layout, PG file conversion, and pattern generation. At the current value of \(\lambda = 3\) microns, blowbacks at approximately 100 to 150 times actual chip dimensions have sufficient detail to enable visual checking of the smallest features. Blowbacks of reticles may also be obtained in the form of color transparencies, to enable
inspection of superposed overlays of various layers.

Once the 10x reticles have been generated, a 1x master mask is made from each reticle using a photorepeater, often called a step and repeat camera. The photorepeater exposes a photographic plate held on a moveable stage, as in the pattern generator. In this case, however, each plate exposure is a 10:1 photo reduction of the reticle pattern. Between exposures the stage is moved by a precise x,y stepping distance. This process is repeated until a complete array of 1x chip patterns for one layer of the system has been exposed. The plate is then developed to produce a 1x master mask. Figure 6 sketches such a mask made from the reticle in figure 5.

Note that when each reticle is inserted in the photorepeater, the position and angular orientation of the reticle pattern is carefully adjusted by microscopic examination of two fiducial marks on the reticle. These marks are placed as part of the pattern generation process, and have the same precise position relative to the chip pattern origin on each of the system's reticles, thus assuring that all mask levels produced with the photorepeater will accurately register with each other.

A succession of contact prints is made from each master mask to yield a number of working masks, sometimes called working plates, for each system layer. These are the actual masks used in wafer fabrication. During the contact printing step of the typical wafer fabrication procedure, the working plates sometimes become worn or damaged, so several are usually made for each layer.

The wafer fabrication facility uses the working plates in the sequence of patterning and process steps described in chapter 2, to produce finished wafers. The fab line requires no detailed information about the design or mask patterns of the integrated system being fabricated. However, several auxiliary patterns are normally included in the mask patterns, some of which are replicated on each chip and examined during wafer fabrication: (i) alignment marks, which are used to accurately overlay successive masks with previous patterning steps, (ii) line width testers, sometimes called critical dimensions (C/D's), which are lines in each mask layer of stated width that may be examined during maskmaking and fabrication to control dimensional tolerances, and (iii) a few simple test transistors and their associated probe pads, which may be electrically tested prior to packaging to verify that the wafer fabrication process was successful.

The finished wafers are divided into chips and packaged by the sequence of steps sketched in figure 7. The wafers are diced into individual chips by first scribing their surface along the
Fig. 2. The **Procedures** and **Artifacts** of the Typical Present Implementation Process
Fig. 3. Sketch Illustrating Function of Pattern Generator

Fig. 4. Parameters of One Flash of Pattern Generator
**Fig. 5. Sketch of a 10x Reticle**

**Fig. 6. Sketch of a Mask made from the above Reticle**
Fig. 7. Sketch Illustrating the Packaging Sequence

(wafer press)
boundary lines between chips, called *scribe lines*, with either a diamond tipped scribe or a diamond edged saw blade, and then fracturing them along these lines. Each individual chip is then cemented into the cavity of a package. After fine wires are bonded between the contact pads on the chip and the leads of the package, and a cover cemented over the cavity, the system is ready for testing.

From the preceding we see that once a system's design files have been produced, all the remaining implementation procedures are design and layout independent, and largely automatic. However, the many extraneous parameters, patterns, and constraints involved in maskmaking and fabrication must be carefully thought through and defined in order to guarantee successful implementation within a reasonable turnaround time. The PG files must be correctly sorted and formatted for the pattern generator to be used. The 10x pattern of the chip must fit within the largest reticle that the pattern generator can produce. The photorepeater used will determine the shape, size, and location of the fiducial marks on the reticle. The size, surface material, and photographic polarity, either positive (*i.e.* clear background field, with opaque features) or negative (*dark field, with transparent features*), of the working plates will be a function of the fabrication facility to be used. Each fab line also typically prescribes its own patterns for the alignment marks and test transistors to be included along with the system in the mask patterns.

While many designs may be scalable and have some longevity, the parameters, patterns, and constraints of maskmaking and fabrication are changing rapidly as the technologies evolve. This constant change complicates interactions with mask and fab firms. Later we describe procedures for implementing moderate sized LSI systems as part of multi-project chips. Such chips are collaborative efforts of many designers, enabling many projects to be merged into one maskmaking and wafer fabrication run. In this way the procedural overhead involved may be shared.
Hand Layout and Digitization using a Symbolic Layout Language

A simple and common method of producing system layouts is to draw them by hand. This is typically done on a one lambda grid using the familiar color codes to identify various system layers. Once the layout has been hand drawn it can then be *digitized*, or translated into machine readable form, by encoding it into a symbolic layout language. Hand layout and digitization using a symbolic layout language is quite a practical method of generating design files for highly structured system designs. Be warned, however, that implementing irregular structures using these primitive procedures is a difficult and tedious task.

If a system has only a few cell types which are replicated over and over, and otherwise has little "random wiring", one need draw only a single copy of each cell type, and then make reproductions or equivalent sized outlines of these cell drawings. All these cell reproductions may then be patched together to plan and build up the overall layout. Similarly, only one symbolic digitization need be made for each cell type. The replication of cells in various orientations and locations in the system layout can then be easily described using the symbolic layout language. In a sense, the ease with which a system's layout can be described using a primitive layout language provides a measure of the regularity of its design. The OM2 Data Chip pictured in the frontispiece was laid out and digitized in this way, using only the simplest machine aids.

The function of a symbolic layout language, in its simplest form, is similar to that of a macro-assembler. The user defines *symbols* (macros) which describe the layout of basic system cells. The locations and orientations of instances of these symbols are described in the language, as a function of appropriate parameters. These symbolic descriptions may then be mechanically processed in a manner similar to the expansion of a macro assembly language program, to yield the intermediate form description of the system layout, which is analogous to machine code, for generating output files. An example intermediate form is described in a later section. The intermediate form files may be processed to yield the PG files: each layer being a machine encoded collection of rectangles encoded as \([x, y, h, w, a]\) values. The generation of PG files is analogous to the loading and execution of machine code to produce output files: it is a process of "unrolling" and fully instantiating all symbol descriptions into a sequence and format suitable for a particular output device. Definition of simple layout languages and the construction of their assemblers is fairly straightforward. The reader may define and implement layout languages by using the macro assembler or higher level language facilities of any commonly available computer system (R1, R3).
The following example will clarify the concepts and procedures of hand layout and symbolic layout description: We wish to create an array of shift registers consisting of parallel horizontal rows of inverters coupled by clocked pass transistors, as in figure 5a., chapter 3. Figure 8a sketches the stick diagram of one row of the array. The entire array can be constructed from one basic cell containing an inverter, the pass transistor following it, VDD and GND buses crossing through on metal, and a clock line passing through on poly. Figure 8b shows a hand sketch of the layout of the basic shift register cell, SRCELL, on a 1\(\lambda\) grid, subject to the design rules given in Ch.2, Sect.2. Since the inverters are coupled by pass transistors, the inverter pullup/pulldown ratio is \(\sim 8:1\) (see Ch.1., Sect.2). Also, while the 4\(\lambda\) wide metal lines could be 1\(\lambda\) narrower in between the contact regions, this would not decrease the cell size. As an exercise, the reader might check for design rule violations, and also for ways of further shrinking the cell size.

The SRCELL layout shown in figure 8b is composed using only rectangles placed at orientations which are integer multiples of 90°. The illustrations and descriptions in this section are considerably simplified by the use of such constrained layout constructions, and yet still illustrate the general principles involved. Were completely arbitrary shapes used, the SRCELL could be made somewhat smaller and still satisfy the design rules. Interestingly, experience has shown that the simple extension of including rectangles at orientations which are integer multiples of 45° enables most cell layouts to reach within a few percent of the minimum area achievable using arbitrary shapes. There is a clear tradeoff here: the inclusion of increasingly complex geometrical objects in a layout will tend to reduce the minimum achievable layout area, but will also increase the computational complexity of the associated machine aids.

We can informally characterize a simple layout language by examining figure 9, which contains a description of the layout of an array of SRCELLs using such a language. The language describes layouts as collections of BOXes on various layers. BOX statements describe each of these boxes by specifying their layer, the X,Y coordinates of their lower left corner and then their length, LX, in the x direction, and LY, in the y direction. The use of a box corner to encode its location simplifies the encoding task. BOX statements may describe arrays of identical boxes, with the array's lower left corner origin at X,Y, by including optional parameters which specify the number NX and replication interval IX in the x direction, and NY and IY in the y direction. Dimensions are given in the length unit, \(\lambda\). A SCALE statement defines the value of \(\lambda\) for this particular layout as \(\lambda = 3.0\) microns.

In figure 9, the SRCELL is first described as a macro, or SYMBOL. The reader can verify that
the collection of BOXes in the definition of the SYMBOL SRCCELL, when ORed together, produces the layout in figure 8b. This SRCCELL is then replicated a number of times in various layout locations according to parameters in several DRAW statements.

Each DRAW statement describes the placement of an array of cells as follows: The cell described by the named SYMBOL definition is considered to be drawn at the origin. It is then mirrored (about the x and/or y axis), and/or rotated (by 0°, 90°, 180°, or 270°) about the origin, as specified by MIRROR or ANGLE transformations. The cell thus positioned may then be replicated NX times at distance intervals IX in the x direction, and that row of cells may then be replicated NY times at intervals IY in the y direction. The resulting array of cells is then translated a distance X,Y from the origin, and placed into the layout.

SCALE LAMBDA = 3.0MICRON:

SYMBOL START, SRCCELL:
BOX DIFF,X = 3,Y = 0,LX = 4,LY = 4,NY = 2,IY = 19;
BOX DIFF,X = 2,Y = 3,LX = 6,LY = 9;
BOX DIFF,X = 8,Y = 8,LX = 3,LY = 2;
BOX DIFF,X = 9,Y = 10,LX = 2,LY = 1;
BOX DIFF,X = 9,Y = 11,LX = 7,LY = 2;
BOX DIFF,X = 16,Y = 9,LX = 4,LY = 4;
BOX DIFF,X = 4,Y = 12,LX = 2,LY = 7;
BOX INV,PL,X = 2.5,Y = 9.5,LX = 5,LY = 10;
BOX POLY,X = 0,Y = 5,LX = 10,LY = 2;
BOX POLY,X = 12,Y = 0,LX = 2,LY = 2b;
BOX POLY,X = 16,Y = 5,LX = 5,LY = 2;
BOX POLY,X = 16,Y = 7,LX = 4,LY = 3;
BOX POLY,X = 2,Y = 11,LX = 6,LY = 7;
BOX CUTS,X = 4,Y = 1,LX = 2,LY = 2,NY = 2,IY = 19;
BOX CUTS,X = 17,Y = 8,LX = 2,LY = 4;
BOX CUTS,X = 4,Y = 9.1,LX = 2,LY = 4;
BOX METL,X = 0,Y = 0,LX = 21,LY = 4,NY = 2,IY = 19; VDD & GND
BOX METL,X = 3,Y = 8,LX = 4,LY = 6;
BOX METL,X = 16,Y = 7,LX = 4,LY = 6;
SYMBOL END;

DRAW SRCCELL,NX = 4,NY = 2,IX = 21,IY = 38,X = 0,Y = 0;
DRAW SRCCELL,MIRRORX,NX = 4,IX = 21,X = 0,Y = 42;

END;

Figure 9. Symbolic Description of Shift Register Array

The "program" in figure 9 describes an array of 3 rows and 4 columns of SRCCELLs. After machine assembly of this program, the resulting design file can be used to generate check plots, which may be inspected to detect errors made in encoding the layout. A check plot of one SRCCELL is given in figure 10a, and we see that the cell has been correctly digitized. A set of
Fig. 8a. Stick Diagram of One Row of Shift Register Array

Fig. 8b. Hand Sketch of Layout of One Shift Register Cell
Figure 10a. Check Plot of the SRCELL

[Dimensions in lambda. Implant layer not shown]

Figure 10b. Check Plot of Stipple Codes
Figure 11. Check Plot of 3 by 4 Array of SRCELLS

[Dimensions in lambda, with the cell outlines indicating relative cell placements according to the program in fig.9]
stipple patterns is used in this check plot to encode the different system layers, with the coding specified in figure 10b. If available, color checkplots are much better: color checkplots can be made denser and still be readable, and association of colors with layers and functions is more easily made and subject to fewer errors in practice. Note: the implant layer hasn't been plotted in fig. 10a, so that the other layers may be more easily seen.

A check plot of the complete 3 by 4 array of cells is given in figure 11 (again the implant layer is not plotted). Although figure 11 is of insufficient scale to check details within the cells, it enables us to check for correct relative placement of the SRCells. The individual cell outlines are included in figure 11, to indicate the nature of the placement of the central row of the array. By mirroring the central row prior to its placement, that row is able to share VDD and GND with the other two rows, thus reducing the overall array size. There is one column of cells per 21 lambda in the x-direction, and one row of cells per 19 lambda in the y-direction. It is very important to note that the outcome of each DRAW statement is determined by the order in which any mirror, rotate, replicate, and translate operations occur (see the section on the Caltech Intermediate Form, and also R2, Ch.6). Any permutation in the order of these operations may lead to a completely different result.

In chapter 3 we found that the PLA is a useful subsystem structure, often used to implement finite state machines and combinatorial logic. We now present a worked out example of a PLA’s layout, to further clarify symbolic layout description. Chapter 3 contains several stick diagrams of PLAs (figs. 13c, 15f). An examination of these stick diagrams reveals that the PLA can be constructed using 6 basic cell types and a slight amount of “random wiring”. Once these 6 basic cells have been laid out by hand and symbolically digitized, it is easy to construct symbolic descriptions of different sized PLAs having various numbers of inputs, product-terms, and outputs.

The digitized layouts of four of these basic cells are check plotted in figure 12. The AND and OR planes of the PLA are constructed as arrays of the 14λ by 14λ PLACellpair cell plotted in figure 12a, which contains two poly and two metal signal lines, and one ground line on the diffusion layer. Diffusion paths may be added in any of four locations in such cells to form transistors, and thus program the PLA. The connection between the AND and OR planes is made using the PLAConnect cell plotted in figure 12b: these cells change the signal paths from the metal to the poly layer. The pullup transistors to be placed at the edges of the AND and OR planes are implemented by the PullupPair cell in figure 12c. The ground return paths, to be
connected to the diffusion lines crossing the planes, are implemented by the PLAground cell in figure 12d. The PLAground cell is structured so that rows of the cell may be inserted at intervals within AND planes, and columns of the cell inserted at intervals within OR planes, to provide proper ground returns in large PLA's. The two other cell types required are the input drivers and output inverters: these cell layouts are left as exercises for the reader. The cells in figure 12 have been collectively planned so as to fit on a 14λ pitch surrounding the PLA's planes. Figure 13 contains a symbolic description of each of these cell types, and a description of a moderate sized PLA constructed from these cells:

Figure 13. Symbolic Description of a 5-Input, 10-Pterm, 8-Output PLA

```
SCALE LAMBDA = 3.0MICRON;

PLA CELL DEFINITIONS:

SYMBOL START,PLACE,PAIR:
BOX DIFF.X = 0,Y = 1,LX = 4,LY = 4,NX = 2,NY = 2; [SEE FIGURE 12A.]

BOX DIFF.X = 8,Y = 0,LX = 2,LY = 14;
BOX POLY.X = 5,Y = 0,LX = 2,LY = 14, NX = 2,IX = 6;
BOX CUTS.X = 1,Y = 2,LX = 2,LY = 2,NY = 2,IY = 7;
BOX METL.X = 0,Y = 1,LX = 14,LY = 4,NY = 2,IY = 7;
SYMBOL END;

SYMBOL START,PLA,CONNECT:
BOX DIFF.X = 0,Y = 1,LX = 4,LY = 4,NY = 2,IY = 7; [SEE FIGURE 12B.]
BOX DIFF.X = 9,Y = 4,LX = 4,LY = 4;
BOX DIFF.X = 13,Y = 4,LX = 3,LY = 2;
BOX POLY.X = 6,Y = 1,LX = 10,LY = 2,NX = 2,IY = 8;
BOX POLY.X = 3,Y = 1,LX = 1,LY = 4,NY = 2,IY = 7;
BOX POLY.X = 14,Y = 7,LX = 2,LY = 2;
BOX CUTS.X = 1,Y = 2,LX = 4,LY = 2,NX = 2,IY = 7;
BOX CUTS.X = 10,Y = 5,LX = 2,LY = 2;
BOX METL.X = 9,Y = 0,LX = 4,LY = 14;
BOX METL.X = 0,Y = 1,LX = 6,LY = 4,NY = 2,IY = 7;
SYMBOL END;

SYMBOL START, PULLUP,PAIR: [SEE FIGURE 12C.]
BOX IMPL.X = 8.5,Y = 0.5,LX = 11,LY = 5;
BOX IMPL.X = 0.5,Y = 4.5,LX = 5,LY = 8;
BOX IMPL.X = 0.5,Y = 7.5,LX = 9,LY = 5;
BOX DIFF.X = 0,Y = 1,LX = 4,LY = 4;
BOX DIFF.X = 4,Y = 2,IX = 16,LY = 2;
BOX DIFF.X = 4,Y = 2,LX = 2,LY = 4;
BOX DIFF.X = 2,Y = 9,IX = 18,LY = 2;
BOX DIFF.X = 9,Y = 8,IX = 4,LY = 4;
BOX POLY.X = 10,Y = 0,IX = 8,LY = 6;
BOX POLY.X = 18,Y = 1,IX = 2,LY = 4;
BOX POLY.X = 8,Y = 8,IX = 2,LY = 4;
BOX POLY.X = 0,Y = 7,IX = 8,LY = 6;
BOX POLY.X = 0,Y = 6,IX = 6,LY = 1;
BOX CUTS.X = 1,Y = 2,LX = 2,LY = 2,NX = 2,IX = 17;
```

Fig. 12a. PLAcellpair

Fig. 12b. PLAconnect

Fig. 12c. PullupPair

Fig. 12d. PLAground

[all dimensions in lambda]
A check plot of the PLA described above is given in figure 14. This check plot has been simplified to include only the outlines of the basic cells, plus the additional wiring necessary to complete the PLA. The dimensions and orientations of the cells may be found by comparing these outlines with the cell details in figure 12. Note that in figure 12 some of the connection points, where paths leave or enter at cell edges or where internal connections may be later inserted, are tagged with tick marks. Cell placements and orientations in the check plot may be visualized by locating and identifying the appropriate connection point marks. A comparison of the check plot with the symbolic description above will clarify the function of the various DRAW statements. To assist in this comparison, the origin cell of the array of cells produced by each DRAW statement has been marked in figure 14 with its cell name. Note that this PLA layout could contain the PLA example of chapter 3, fig. 15f.

Symbolic layout languages are easy to define, and may be primitive or sophisticated, according to the requirements of the user. The function of the assembler for such a language is simply to scan and decode the statements and translate them into design files in intermediate form. Conversion of design files into check plot or pattern generator output files is straightforward for the above simple language, since we have used only boxes with a severe constraint on angular orientations. MIRROR and ANGLE transformations are easily handled: x and y coordinates of symbols and boxes are simply replaced by ±x or ±y, according to the specific parameters, during the instantiation of symbols and drawing of boxes prior to their replication and translation into the layout output file.

The effectiveness of the above language could be further increased by constructing an assembler capable of handling nested symbols. By using nested symbols, system layouts may be described in a hierarchical manner, leading to very compact descriptions of structured designs. At the lowest level, one might define symbols for such small but commonly encountered structures as the various forms of contacts. Boxes and these simple symbols could then be used to construct cells such as those in the PLA example above. The PLA could be constructed with these cells, and then defined as a symbol to be used in a larger design. An example of the sort of function
Fig. 14. Check Plot using only Cell Outlines, of the 5-Input, 10-Pterm, 8-Output PLA

(dimensions in lambdas; symbol labels on origin cells of the fig.13 DRAW statements)
Fig. 15. Design, Hand Layout, Design File Generation, And Design Checking Using a Symbolic Layout Language and Simple Machine Aids
one might add to create a much more sophisticated language, and language processor, would be the capability of generating the layout description of a PLA from the collection of basic cells, as a function of its input, output, and output size parameters and logic function parameters.

Figure 15 summarizes the procedures and artifacts of hand layout, and layout description and digitization using a layout language. By studying figure 15, and thinking back over the material and examples of this section, one can visualize a complete, though primitive, sequence of steps sufficient to prepare a design for implementation. These procedures are entirely adequate for preparing small LSI projects for implementation. The procedures may also be used for those large LSI systems which have highly structured designs.

The primary obstacle that these primitive procedures place in the path of the system designer is the sheer time and effort it takes to get through the loop to a new check plot each time a small design change is made. The enthusiasm aroused by a sudden insight, such as the conception of a completely new topological possibility for an important system cell, can be dampened by the tedious tasks of hand layout and box digitization required before one can really see the full effect of the idea on the overall system layout.

Though often supported by large batch mode CAD systems for containing, modifying, checkplotting, and simulating designs, the majority of LSI layout now done in industry begins with hand layout. Digitization is usually simplified by the use of digitizing tables, which are much like graphics plotters in reverse: a new section of a design, laid out by hand, is placed on the table and digitized by tapping switches while manually following the outlines of the cell's boxes with a pointer. Although this is less tedious than digitization using a layout language, it is still time-consuming and hardly interactive.

The next section describes an interactive graphics layout system which enables the system designer to quickly sketch new layout ideas and see their effect immediately.
An Interactive Layout System

[Section contributed by Douglas Fairbairn, Xerox PARC, and James Rowson, Caltech]

Computing hardware of sufficient power to support highly interactive graphics has in the past been quite expensive, and this has inhibited the widespread application of interactive computing techniques. However, because of expected advances in VLSI technology, we are rapidly approaching the day when many will have access to personal computers with computing power rivaling today's medium to large-scale systems. It will be more difficult to provide effective software for these systems than it will be to build the computers themselves.\(^2\) In this section we describe a highly interactive layout system which runs on a modest personal computer, rather than on an expensive, limited access, centralized system. This system was developed anticipating the work environment of the future, in which most "knowledge" workers will have personal computers as part of their normal office equipment.

ICARUS\(^1\) (Integrated Circuit ARtwork Utility System) is a software system which enables the user to create and modify an integrated system layout directly on a CRT display screen. ICARUS was conceived with the idea that the designer would create and edit a layout at the display, without doing any more than a rough sketch or "stick diagram" before beginning work. Creating and moving items is fast and easy enough so that the designer can truly sketch on the screen. Once the layout is basically correct, the items can be moved or modified to arrive at the most compact layout.

The user is required to remember very little about the available commands or their use because the commands themselves are displayed on the screen and the system prompts the user for additional information as it is needed. The system can format and output check plots to matrix type printers or on raster-scan laser printers. ICARUS design files can be used to create standard pattern generation files from which masks can be made. An overview of design and layout procedures using the system is given in figure 16. It is instructive to compare this with figure 15, which presents equivalent steps for hand layout.

All the software to accomplish these various steps runs on a small experimental minicomputer known as the Alto. This machine was designed by researchers at Xerox PARC as a general purpose personal computer suitable for both text and graphics applications\(^2\). No additional, special hardware is used by ICARUS. The ICARUS system is programmed in BCPL, an ALGOL-like high level language. There are about 30K words of compiled code in the system of
Fig. 16. Design, Layout, Design File Generation, and Design Checking Using an Interactive Graphics Layout System
which half is in memory at any given time. At minimum, the Alto memory has 64K 16 bit words. A 2.5 Mbyte cartridge disk drive is an integral part of the system. The user interacts with the system through an encoded keyboard (software definable keys) and with a pointing device called a mouse (R2, p173). A cursor is controlled on the screen by moving the mouse around on a small area of the user’s desk. A bit map display with a resolution of 600x800 dots is used for output, and printers for doing check plots are available through an in-house computer network.

The ICARUS display features two windows which provide a flexible working view of the layout, as shown in figure 17. The upper window is normally used for viewing a large piece of the layout at small magnification, and the lower window used for looking at a smaller section in more detail. The magnifications of the windows may be set independently.

In addition to the two windows there are various menus and status lines presented in the display. The menu on the left is the command menu. The menu under the upper window is the parameter menu. Under the parameter menu is the stipple menu, containing the mask level codes. Rectangles at a given level are stippled with the pattern for that level. The patterns were chosen so that, where necessary, one pattern could be seen through the other to verify that appropriate layers are overlapping properly. Current drawing coordinates and the status of system memory space are displayed to the right of the stipple menu.

The user interface is implemented principally through the display, the mouse and five conveniently located keys on the keyboard. Frequently used commands are given using only one or two simple hand operations, and can be done without glancing away from the display. These characteristics, coupled with rapid display redrawing, enhance the system’s interactiveness.

The internal data representation in ICARUS is based on three types of items: rectangles, symbols, and text strings. The organization of these items into memory data structures, and the typical run-time memory space allocation is illustrated in figure 18.

Rectangles are created with the aid of the mouse. They may have angular orientations which are integer multiples of 45°. They can be moved, copied, or deleted using the mouse and one key. As items are created, they are added to an item list in main memory. Each rectangle is stored as 6 words in memory: the first word is the pointer to the next item, the second specifies what layer it is on, what type of item it is, etc. The third through sixth words specify the minimum and maximum x and y coordinates. The items are kept in order of increasing values of minimum x.
coordinate, so that the display may be quickly redrawn.

When a symbol is defined by the user, the items which are contained within it are stored on the disk, while a pointer, the name and the bounding box for the symbol are placed in main memory. Symbols can be nested to any level. Once a symbol definition has been created, one is free to define symbol instances which are references to that definition. The symbol instance may be a command to draw one copy of the symbol at a certain location, or a whole array. The size of the symbol instance, which resides in main memory, is the same in both cases. The use of symbols wherever possible tends to preserve main memory space. Rather large systems can be designed using ICARUS, if the systems are well structured and make extensive use of symbols. This is true even when using a minimum sized 64K memory, which leaves little space for layout data.

Text is used for identifying data and control lines and is merely a memory aid to the user. There is no attempt to make use of the text or other information in the drawing for connectivity or other types of checking.

Operations more complex than those such as draw and move are implemented through the use of menus as shown in figure 17. The desired command is chosen by pointing at it with the cursor and clicking a mouse button. The selected command is then inverted to white on black video to identify its selection, which the user then confirms with a key on the keyboard. At this point, the system prompts the user with instructions presented in the display area normally holding the stipple menu. The instructions lead the user through the individual steps required, for example, to mirror or rotate a group of items.

Operations on symbols are defined in a secondary menu which can be reached by selecting the command "symbols" on the primary menu. The secondary menu offers commands such as define symbol, draw symbol, list the names of the symbols in the symbol library, or expand symbol. This last command is used to modify a symbol which is already defined, the modified symbol definition immediately updating all symbol instances which point to it.

Various system parameters are displayed in the parameter line directly below the top window. Values such as the default line width for the currently selected layer, the magnification of the top and bottom windows, and the spacing of the tick marks are all displayed. The parameter values can be changed at any time by selecting the desired one and typing the new parameter value on the keyboard. The X,Y layout coordinates of the point last clicked with the mouse are displayed.
Fig. 17. The ICARUS Display: 2 Views of a Layout in Progress
Figure 18. ICARUS Memory Allocation and Data Structure
at the right of the screen. The DX, DY distances between the last two clicks are also displayed. This feature provides a convenient "ruler" for measuring distances on the layout.

The construction of an interactive layout system such as ICARUS is a relatively straightforward task for one who is experienced in interactive computer graphics (R2), given a display oriented minicomputer system and effective systems building software. A first version of ICARUS was constructed in 3 man-months, and a mature version produced in an additional 5 man-months.

ICARUS has been used internally in Xerox to lay out many integrated system projects, and to organize a number of multi project chips. Among the users were a number of individuals previously unfamiliar with integrated circuit layout, who nevertheless successfully completed LSI projects with up to 10,000 transistors. We find that the interactive nature of such a system not only aids the experienced designer, but also enhances the learning process for the novice. We believe that such interactive, personal design systems greatly enhance the creative ability of the designer by enabling easy generation and examination of many more design alternatives per unit time than would be the case with centralized, non-interactive design systems.

However, there is more to integrated system design than circuit layout. Design rules must be checked, logic transfer functions tested, and, in certain cases, circuit transfer functions computed to determine delays and predict system performance. We believe that the direction in which to search for further improvements in design tools is in the replacement of the primitive ICARUS type of data structure with one which allows design functions other than just layout to also interactively operate upon the same data base. This is the subject of the later section on fully integrated design systems.
The Caltech Intermediate Form for LSI Layout Description

[Section Contributed by Robert F. Sproull, Carnegie-Mellon University, and Richard F. Lyon, Xerox PARC]

The Caltech Intermediate Form (CIF Version 2.0) is a means of describing graphic items (mask features) of interest to LSI circuit and system designers. Its purpose is to serve as a standard machine readable representation from which other forms can be constructed for specific output devices such as plotters, video displays, and pattern-generation machines. The intermediate form is not intended as a symbolic layout language: CIF files will usually be created by computer programs from other representations, such as a symbolic layout language or an interactive design program. Nevertheless, the form is a fairly readable text file, in order to simplify combining files and tracing difficulties.

The basic idea of the form is to specify literally every geometric object in the design using ample precision. Use of this form provides participating design groups easy access to output devices other than their own, enables sharing designs with others, allows combining several designs to form a larger chip, and the like. It is not necessary for all participating groups to implement the entire set of features of CIF, as long as their programs and documents contain warnings about unimplemented functions; nevertheless, the syntax must be correctly interpreted by all programs that read CIF, to assure a reasonable result.

CIF thus serves as the common denominator in the descriptions of various integrated system projects. No matter what the original input methods are (hand layout and coding, or a design system), the designs will be translated to CIF as an intermediate, before being translated again to a variety of formats for output devices or other design aids.

The original CIF was conceived by Ivan Sutherland and Ron Ayers in 1976. Subsequent improvements were contributed by Carlo Sequin, Douglas Fairbairn, and Stephen Trimberger.

This specification is divided into four parts: a description of the syntax of the form, a description of the semantics, an explanation of the transformations used, and a discussion of the conversion of wires to boxes.
Syntax

A CIF file is composed of a sequence of characters in a limited character set. The file contains a list of commands, followed by an end marker; the commands are separated with semicolons. Commands are:

<table>
<thead>
<tr>
<th>Command</th>
<th>Form</th>
</tr>
</thead>
<tbody>
<tr>
<td>Polygon with a path</td>
<td>P path</td>
</tr>
<tr>
<td>Box with length, width, center, and direction (direction defaults to (1,0) if omitted)</td>
<td>B integer integer point point</td>
</tr>
<tr>
<td>Round flash with diameter and center</td>
<td>R integer point</td>
</tr>
<tr>
<td>Wire with width and path</td>
<td>W integer path</td>
</tr>
<tr>
<td>Layer specification</td>
<td>L shortname</td>
</tr>
<tr>
<td>Start symbol definition with index, a, b (a and b both default to 1 if omitted)</td>
<td>DS integer integer integer</td>
</tr>
<tr>
<td>Finish symbol definition</td>
<td>DF</td>
</tr>
<tr>
<td>Delete symbol definitions</td>
<td>DD integer</td>
</tr>
<tr>
<td>Call symbol</td>
<td>C integer transformation</td>
</tr>
<tr>
<td>User extension</td>
<td>digit userText</td>
</tr>
<tr>
<td>Comments with arbitrary text</td>
<td>( commentText )</td>
</tr>
<tr>
<td>End marker</td>
<td>E</td>
</tr>
</tbody>
</table>

A more formal definition of the syntax is given below. The standard notation proposed by Niklaus Wirth\(^{14}\) is used: production rules use equals = to relate identifiers to expressions, vertical bar | for or, and double quotes " " around terminal characters; curly brackets { } indicate repetition any number of times including zero; square brackets [ ] indicate optional factors (i.e. zero or one repetition); parentheses ( ) are used for grouping; rules are terminated by period. Note that the syntax allows blanks before and after commands, and blanks or other kinds of separators (almost any character) before integers, etc. The syntax reflects the fact that symbol definitions may not nest.

cifFile
command = { { blank } { command } semi } endCommand { blank }.
primCommand = primCommand | defDeleteCommand |
defStartCommand semi { { blank } { primCommand } semi } defFinishCommand.
polygonCommand | boxCommand | roundFlashCommand | wireCommand | layerCommand | callCommand | userExtensionCommand | commentCommand.

polygonCommand = "P" path.
boxCommand = "B" integer sep integer sep point [ sep point ].
roundFlashCommand = "R" integer sep point.
wireCommand = "W" integer sep path.
layerCommand = "L" { blank } shortname.
defStartCommand = "D" { blank } "S" integer [ sep integer sep integer ].
defFinishCommand = "D" { blank } "F".
defDeleteCommand = "D" { blank } "D" integer.
callCommand = "C" integer transformation.
userExtensionCommand = digit userText.
commentCommand = "(" commentText "),".
endCommand = "E".

transformation = { blank } { "T" point | "M" { blank } "X" | "M" { blank } "Y" | "R" point }.
path = point { sep point }.
point = sinteger sep sinteger.

integer = { sep } [ "-" ] integerD.
integerD = { sep } integerD.
digit = digit { digit }.

shortname = c [ c ] { c } [ c ].
c = digit | upperChar.
userText = { userChar }.
commentText = commentChar | commentText "(" commentText ")" commentText.

semi = { blank } ";" { blank }.
sep = upperChar | blank.
digit = "0" | "1" | "2" | "3" | "4" | "5" | "6" | "7" | "8" | "9".
upperChar = "A" | "B" | "C" | ... | "Z".
blank = any ASCII character except digit, upperChar, ";", "(" or ")".
userChar = any ASCII character except ";".
commentChar = any ASCII character except "(" or ")".

Semantics

The fundamental idea of the intermediate form is to describe unambiguously the geometry of patterns for LSI circuits and systems. Consequently, it is important that all readers and writers of files in this form have exactly the same understanding of how the file is to be interpreted. Many of the decisions in designing the file format were made to avoid ambiguity or small but troublesome errors: floating point numbers are avoided; there are no iterative constructs, though there may be in future additions to CIF.

A simple file format might include only primitive geometric constructs, such as polygons, boxes, flashes and wires. Unfortunately, the geometric description of a chip with hundreds of thousands of rectangles on it would require an immense file of this sort. Consequently, we have made provision for defining and calling symbols; this should reduce the size of the file substantially.

It is important that programs processing CIF files operate cautiously, maintaining a constant vigilance for mistakes or entries that will not be processed properly. The description below mentions implementation suggestions or cause for caution inside brackets [ ].
Measurements. The intermediate form uses a right-handed coordinate system shown in Figure 19a, with x increasing to the right and y increasing upward. (Directions and distances are always interpreted in terms of the front surface of the finished chip, not in terms of the various sizes and mirrorings of the intermediate artifacts.) The units of distance measurement are hundredths of a micron (μm); there is no limit on the size of a number. [Programs reading numbers from CIF files should check carefully to be sure that the number does not overflow the number of bits in the internal representation used, and should specify their own limits, if any.]

Directions. Rather than measure rotation by angles, CIF uses a pair of integers to specify a "direction vector." (This eliminates the need for trigonometric functions in many applications, and avoids the problem of choosing units of angular measure.) The first integer is the component of the direction vector along the x axis; the second integer along the y axis. Thus a direction vector pointing to the right (the +x axis) could be represented as direction (1 0), or equivalently as direction (17 0); in fact, the first number can be any positive integer as long as the second is zero. A direction vector pointing NorthEast (i.e., rotated 45 degrees counterclockwise from the x axis) would have direction (1 1), or equivalently (3 3), and so on. [A (0 0) direction vector may be defaulted to mean the +x axis; a warning should be generated.]

Geometric primitives. The various primitives that specify geometric objects are not intended to be mutually exclusive or exhaustive. CIF may be extended occasionally to accommodate more exotic geometries. At the same time, it is not necessary to use a primitive just because it is provided. Notice in the examples below that lower case comments and other characters within a command are treated as blanks, and that blanks and upper case characters are acceptable separators.

Boxes: Box Width 60 Length 25 Center 80,40 Direction -20,20; (or B60 25 80 40 -20 20;) The fields which define a box are shown graphically in Figure 19a. Center and direction (optional, defaults to +x axis) specify the position and orientation of the box, respectively. Length is the dimension of the box parallel to the direction, and Width is the dimension perpendicular to the direction.

Polygons: Polygon A 0.0 B 10.20 C -30.40; (or P0 0 10 20 -30 40;) A polygon is an enclosed region determined by the vertices given in the path, in order. For a polygon with n sides, n vertices are specified in the path (the edge connecting the last vertex with the first is implied; see Figure 19b). [Programs that try to interpret polygons may place various restrictions on
their paths; no set of constraints has been generally accepted, and no program currently exists for converting completely general polygons to pattern generator output.

**Flashes**: RoundFlash Diam 200 Center 0-00,000; (or R200 0-00,000;)

The diameter of a flash is sufficient to specify its shape, and the center specifies its position. (see Figure 19b). [Some programs may substitute octagons, or other approximations, for round flashes.]

**Wires**: Wire Width 50 A 0,0 B 10,20 C -30,40; (or W50 0 0 10 20 -30 40;)

It is sometimes convenient to describe a long, uniform width run by the path along its centerline. We call this construct a wire (see Figure 19b). An ideal wire is the locus of points within one half-width of the given path. Each segment of the ideal wire therefore includes semicircular caps on both ends. Connecting segments of the wire is a transparent operation, as is connecting new wires to an existing one: the semicircular overlap ensures a smooth connection between segments in a wire and between touching wires. [For output devices that have a hard time constructing circles, we approximate the ideal wire with squared-off ends. Notice that squared-off ends work nicely for segments meeting at right angles, but cause problems if wires or wire segments are connected at arbitrary angles. A way to circumvent this problem is to convert, prior to output, any wires in a file into connected sets of boxes of appropriate length, width, angle and center position (Figure 19c). The width of each box is the same as the width of the wire. The length of the boxes must be adjusted to minimize unfilled wedges and overlapping "ears". An algorithm for constructing boxes from a wire description is given in a later subsection. If the wire is specified within a symbol definition, the approximation need be computed only once, and can then be used each time the symbol is instantiated.]

**Layer specification**: Layer ND nmos diffusion; (or LND;)

Each primitive geometry element (polygon, box, flash, or wire) must be labeled with the exact name of a fabrication mask on which it belongs. Rather than cite the name of the layer for each primitive separately, the layer is specified as a "mode" that applies to all subsequent primitives, until the layer is set again (layer mode is preserved across symbol calls, which are discussed later).

The argument to the layer specification is a short name of the layer. Names are used to improve the legibility of the file, and to avoid interfering with the various biases of designers and fabricators about numbers (one person's "first layer" is another's "last"). [The intention of the layer specification command is to label locally the layer for a particular geometry. It is therefore senseless to specify a box, wire, polygon or flash if no layer has been specified. In order to detect this error, the command LZZZZ is implicitly inserted at the beginning of the file, and as the first command of a symbol definition (DS: see below). Any attempt to generate geometric output on layer ZZZZ will result in an error.]
It is important that layer names be unique, so that combining several files in intermediate form will not generate conflicts. The general idea is that the first character of the name denotes the technology, and the remainder is mnemonic for the layer. At present, the following layers are defined:

<table>
<thead>
<tr>
<th>Layer</th>
<th>Description</th>
</tr>
</thead>
<tbody>
<tr>
<td>ND</td>
<td>NMOS Diffusion</td>
</tr>
<tr>
<td>NP</td>
<td>NMOS Polysilicon</td>
</tr>
<tr>
<td>NC</td>
<td>NMOS Contact cut</td>
</tr>
<tr>
<td>NM</td>
<td>NMOS Metal</td>
</tr>
<tr>
<td>NI</td>
<td>NMOS depletion mode implant</td>
</tr>
<tr>
<td>NB</td>
<td>NMOS Buried contact</td>
</tr>
<tr>
<td>NG</td>
<td>NMOS overGlass openings</td>
</tr>
</tbody>
</table>

New layer names will be defined as needed.

Programs that read CIF will want to check to be sure that layer names used do in fact correspond to fabrication masks being constructed. However, the file may cite layer names not used in a particular pass over the CIF file. It would be helpful for the program to provide a list of the layer names that it ignored.

Symbols. Because many LSI layouts include items that are often repeated, it is helpful to define often-used items as "symbols." This facility, together with the ability to "call" for an instance of the symbol to be generated at a specific position, greatly reduces the bulk of the intermediate form.

The symbol facilities are deliberately limited, in order to avoid mushrooming difficulties of implementing programs that process CIF files. For example, symbols have no parameters; calling a symbol does not allow the symbol geometry to be scaled up or down; there are no direct facilities for iteration. The main reason for symbol facilities is to limit the file size: if the symbol mechanism is not adequate for some application, the desired geometry can still be achieved with less use of symbols, and more use of explicit geometrical primitives. [Symbols need not be used at all: this eliminates the need for intermediate storage for symbol definitions, but results in larger design files. Machines which must process a fully-instantiated representation of a layer (such as pattern generators) might only accept CIF files without symbol definitions, to reduce the cost of implementation. Therefore, it would be useful to have a program that would convert general CIF files to fully instantiated CIF files, and maybe to sort by layer, location, or whatever.]

The ability to call for iterations (arrays) of symbols is not provided in CIF Version 2.0. This is primarily due to the difficulty of defining a standard method of specifying iterations, without
introducing machine-dependent computation problems. It is still possible to achieve a great deal of file compaction by defining several layers of symbols (e.g. cell, row, double-row, array, etc.). However, the ability to iterate symbol calls is a likely prospect for a future addition to CIF.

**Defining symbols:** Definition Start #57 A/B = 100/1; ... ; Definition Finish; (or DS57 100 1; ... ;DF)
A symbol is defined by preceding the symbol geometry with the DS command, and following it with the DF command. The first argument of the DS command is an identifying symbol number (unrelated to the order of listing of symbol definitions in the file).

The mechanism for symbol definition includes a convenient way to scale distance measurements. The second and third arguments to the DS command are called a and b respectively. As the intermediate form is read, each distance (position or size) measurement cited in the various commands (polygons, boxes, flashes, wires and calls) in the symbol definition is scaled to \((a \times \text{distance})/b\). For example, if the designer uses a grid of 1 micron, the symbol definition might cite all distances in microns, and specify \(a = 100\), \(b = 1\). Or the designer might choose lambda (characteristic fabrication dimension) as a convenient unit. This mechanism reduces the number of characters in the file by shrinking the integers that specify dimensions and may improve the legibility of the file (it does not provide scaling, or the ability to change the size of a symbol called within the definition).

Definitions may not nest. That is, after a DS command is specified, the terminating DF must come before the next DS. The definition may, however, contain calls to other symbols, which may in turn call other symbols.

There is only one restriction on the placement of symbol definitions in the file: a symbol must be defined before its instantiation becomes necessary. This constraint can be satisfied by placing all symbol definitions first in the file, and then calls on the symbols. In fact, it is often convenient to have the file consist exclusively of symbol definitions and ONE call on a symbol. This call will be the last command in the file before the end command. [If a definition redefines a symbol that already exists, the previous definition is discarded: a warning message should be generated. When several people contribute to a design, some symbol management is therefore necessary; see Deleting symbol definitions below.]

**Calling symbols:** Call Symbol #57 Mirrored in X Rotated to \(-1,1\) then Translated to \(10,20\);
The `C` command is used to call a specified symbol and to specify a transformation that should be
applied to all the geometry contained in the symbol definition. The call command identifies the symbol to be called with its "symbol index," established when the symbol was defined.

The transformation to be applied to the symbol is specified by a list of primitive transformations given in the call command. The primitive transformations are:

- **T point**: Translate the current symbol origin to this point.
- **M X**: Mirror in X, i.e., multiply X coordinate by -1.
- **M Y**: Mirror in Y, i.e., multiply Y coordinate by -1.
- **R point**: Rotate symbol's x axis to this direction.

Intuitively, each coordinate given in the symbol is transformed according to the first primitive transformation in the call command, then according to the second, etc. Thus "C1 T500 O MX" will first add 500 to each x coordinate from symbol 1, then multiply the x coordinate by -1. However, "C1 MX T500 O" will first multiply the x coordinate by -1, and then add 500 to it: the order of application of the transformations is therefore important. In order to implement the transformations, it is not necessary to perform each primitive operation separately; the several operations can be combined into one matrix multiplication (see the subsection on transformations).

Symbol calls may nest; that is, a symbol definition may contain a call to another symbol. When calls nest, it is necessary to "concatenate" the effects of the transformations specified in the various calls (see the subsection on transformations). [There is no sensible way in which a symbol may be invoked recursively (i.e., call itself, either directly or indirectly). Programs that read the intermediate form should check that no recursion occurs. This can be achieved by retaining a single flag with each symbol to indicate whether the symbol is currently being instantiated; the flags are initialized to "false." When a symbol is about to be instantiated, we check the flag; if it is "true," we have detected recursion, print an error message and do not perform the call. Otherwise, we mark the flag "true," instantiate the symbol as specified, and mark the flag "false" when the instantiation is complete.]

Layer settings are preserved across symbol calls and definitions. Thus, in the sequence:

```
LMN;
S6 20 0;
C 57 T45 13;
DS 114...;
DF;
LMN;
S3 0 0;
```
the second LNM is not necessary, regardless of the specification of symbols 57 and 114.

 Deleting symbol definitions: Delete Definitions greater than or equal to 100; (or DD100;)
The DD command signals the program reading the file that all symbols with indices greater than
or equal to the argument to DD can be "forgotten" -- they will not be instantiated again. This
feature is included so that several intermediate form files can be appended and processed as one.
In such a case, it is essential to delete symbol definitions used in the first part of the file both
because the definitions may conflict with definitions made later and because a great deal of
storage can usually be saved by discarding the old definitions.

The argument to DD that allows some definitions to be kept and some deleted is intended to be
used in conjunction with a standard "library" of definitions that a group may develop. For
example, suppose we use symbol indices in the range 0 to 99 for standard symbols (pullup
transistors, contacts, etc.) and want to design a chip that has 2 student projects on it. Each project
defines symbols with indices 100 or greater. The CIF file will look like:

```
/Definitions of library symbols;
DS 0 100 1;
/...definition of symbol 0 in library;
DF;
DS 1 100 1;
/...definition of symbol 1;
DF;
/...remainder of library;

/Begin project 1;
DS100 100 1;
/...first student's first symbol definition;
DF;
...
DS109 100 1;
/...first student's main symbol definition;
DF;
C109 T403 -110; / call on first student's main symbol;
DD100; / Preserve only symbols 1 to 99;

/Begin project 2;
DS100 100 1;
/...second student's first symbol definition;
DF;
...
DS113 100 1;
/...second student's main symbol definition;
C1 T-3 45; / Call on library symbol, still available;
DF;
C113 T401 0; / call on second student's main symbol;
E
```
User expansion: 3'SYMBOL.LIBRARY'; 5' NONSTANDARD DESIGN RULES: LAMBDA = 4.0;

Several command formats (any command starting with a digit) are reserved for expansion by individual users; the authors of the intermediate form agree never to use these formats in future expansions of the standard format. For example, private expansions might provide for (1) requesting that another file be "inserted" at this point in the processing, thus simplifying the use of symbol libraries; (2) inserting instructions to a preprocessor that will be ignored by any program reading only standard intermediate form constructs; or (3) recording ancillary information or data structures (e.g., circuit diagrams, design-rule check results) that are to be maintained in parallel with the geometry specified in the style of the intermediate form.

Comments: (HISTORY OF THIS DESIGN);

The comment facility is provided simply to make the file easier to read. [It is possible to deactivate any number of commands by simply enclosing them within a pair of parentheses, even if they already include balanced parentheses.]

End Command: End of file.

The final E signals the end of the CIF file. [Programs that read CIF should give an error message if the file ends without an End Command, or a warning if more text other than blanks follows the E.]
Transformations (see also reference R2)

When we are expanding a symbol, we need to apply a transformation to the specification of an item in the symbol definition to get the specification into the coordinate system of the chip. There are three sorts of measurements that must be transformed: distances (for widths, lengths), absolute coordinates (for "points" in all primitives) and directions (for boxes).

Distances are never changed by a symbol call, because we allow no scaling in the call. Thus a distance requires no transformation.

A point \((x, y)\) given in the symbol is transformed to a point \((x', y')\) in the chip coordinate system by a \(3 \times 3\) transformation matrix \(T\):

\[
[x' \ y' \ 1] = [x \ y \ 1] \ T
\]

[It is a good idea to check either the last column of \(T\), or the 1 at the end of the transformed vector, even though they never need to be computed.]

\(T\) is itself the product of primitive transformations specified in the call: \(T = T_1 T_2 T_3\), where \(T_1\) is a primitive transformation matrix obtained from the first transformation primitive given in the call, \(T_2\) from the second, and \(T_3\) from the third (of course, there may be fewer or more than 3 primitive transformations specified in the call). These matrices are obtained using the following templates for each kind of primitive transformation:

- \(T\) a b.
  \[
  T_n = \begin{bmatrix}
  1 & 0 & 0 \\
  0 & 1 & 0 \\
  a & b & 1
  \end{bmatrix}
  \]

- M X.
  \[
  T_n = \begin{bmatrix}
  1 & 0 & 0 \\
  0 & 1 & 0 \\
  0 & 0 & 1
  \end{bmatrix}
  \]

- M Y.
  \[
  T_n = \begin{bmatrix}
  1 & 0 & 0 \\
  0 & -1 & 0 \\
  0 & 0 & 1
  \end{bmatrix}
  \]

- R a b.
  \[
  T_n = \begin{bmatrix}
  a/c & b/c & 0 \\
  -b/c & a/c & 0 \\
  0 & 0 & 1
  \end{bmatrix}
  \text{ where } c = \text{Sqrt}(a^2 + b^2)
  \]

Transformation of direction vectors \((x \ y)\) is slightly different than the transformation of
coordinates. We form the vector \([x \ y \ 0]\), and transform it by \(T\) into the new vector \([x' \ y' \ 0]\). The transformed direction vector is simply \((x' \ y')\). (Note that some output devices may require rotations to be specified by angles, rather than direction vectors. Conversion into this form may be delayed until necessary to generate the output file. Then we calculate the angle as \(\text{ArcTan}(y/x)\), applying care when \(x=0\).]

Nested calls require that we combine the transformations already in effect with those specified in the new call. Suppose we are expanding a symbol \(a\), as described above, transforming each coordinate in the symbol to a coordinate on the chip by applying matrix \(Tac\). Now we encounter, in \(a\)'s definition, a call to \(b\). What is to happen to coordinates specified in \(b\)? Clearly, the transformations specified in the call will yield a matrix \(Tba\) that will transform coordinates specified in symbol \(b\) to the coordinate system used in symbol \(a\). Now these must be transformed by \(Tac\) to convert from the system of symbol \(a\) to that of the chip. Thus, the full transformation becomes

\[
[x' \ y' \ 1] = [x \ y \ 1] Tba Tac
\]

The two matrices may be multiplied together to form one transformation \(Tbc = (Tba Tac)\) that can be applied to convert directly from the coordinates in symbol \(b\) to the chip. This procedure can be carried to an arbitrary depth of nesting.

To implement transformations, we proceed as follows: we maintain a "current transformation matrix" \(T\), which is initialized to the identity matrix. We use this matrix to transform all coordinates. When we encounter a symbol call, we:

1. "Push" the current transformation and layer name on a stack.
2. Set layer name to \(ZZZZ\).
3. Collect the individual primitive transformations specified in the call into the matrices \(T1, T2, T3\) etc.
4. Replace the current transformation \(T\) with \(T1 T2 T4 \ldots T\); i.e., premultiply the existing transformation by the new primitive transformations, in order.
5. Now process the symbol, using the new \(T\) matrix.
6. When we have completed the symbol expansion, "pop" the saved matrix and layer name from the stack. This restores the transformation to its state immediately before the call.
Decomposing Wires Into Boxes

The following algorithm for decomposing wires into boxes was developed by Carver Mead, and first implemented at Caltech by Ron Ayers; it was further modified to be consistent with the use of direction vectors, to allow more general path lengths, and to avoid use of trigonometric functions. [Note that this decomposition covers more area than the locus of points within w/2 of the path for small angles of bend, but less area for sufficiently sharp bends; in particular, if a path bends by 180 degrees (reverses) it will have no extension past the point of reversal (it is missing a full semicircle). Other decompositions are possible, and may better approximate the correct shape.]

Let the wire consist of a path of n points \( p_1 \ldots p_n \).
Let \( w \) represent the width of the wire.

"Initialization:"
IF \( n = 0 \) THEN DONE; "no path"
IF \( n = 1 \) THEN
\{MAKEFLASH[Diameter = w, Center = \( p_1 \)]; "single-point gets a flash";\}
DONE;

\( i = 1; \)
OldEntension = \( w/2 \); "initial end of wire"
Segment = \( p_2 - p_1 \); "\( p_1 \) and \( p_2 \) are points in path, Segment is a vector (a point)"
"LoopConditions:" 
FOR \( p_i \) to \( p_{i+1} \) in path UNTIL \( p_{i+1} \) is last DO
"calculate the box for the segment from \( p_i \) to \( p_{i+1} \):"
IF \( p_{i+1} \) is last THEN \{ Entension = \( w/2 \); "final end of wire" \} 
ELSE 
{ "compute Entension for intermediate point:" 
NextSegment = \( p_{i+2} - p_{i+1} \); "next vector in path"
T = MATRIX[ 
\[x[Segment], -y[Segment],
\[y[Segment], x[Segment] \];
"T transforms Segment to +x axis."
Bend = MULTIPLE[ NextSegment, T ]; "relative direction vector"
"if Bend is \( (0,0) \), delete \( p_{i+1} \), reduce \( n \), and start over"
Entension = \( w/2 \) * ( ABS[y[Bend]] / 
\( \text{LENGTH}[\text{Bend}] + \text{ABS}[x[\text{Bend}]] \) );
};
MAKEBOX [ 
\{ Length = \text{LENGTH}[\text{Segment}] + \text{Entension} + \text{OldEntension}; \},
\{ Width = w; \},
\{ Center = \( (p_i + p_{i+1})/2 + (\text{Segment} / \text{LENGTH}[\text{Segment}]) \) * 
(\text{Entension} - \text{OldEntension})/2; \},
\{ Direction = \text{Segment}; "careful, may be zero vector" \} ];
\( i = i + 1; \)
OldEntension = \text{Entension};
Segment = NextSegment; "next vector in path"
ENDLOOP;
DONE;
Fig. 19a. Box Representation in Intermediate Form

Fig. 19b. Other Items in the Intermediate Form

T transforms Segment to the +X axis

- AB = Segment * T
- BC = NextSegment * T
- Bend = Vector BC
- Extension = BG = BH

Similar triangles BCD, EFG, BFH

- BC,CD,DB :: EF,FG,GE :: BF,FH,HB
- FG = FB + BG
  - = BH * (BC/DB) + BG
  - = (1 + BC/DB) * BG
- BG = FG / (1 + BC/DB)
  - = GE * (BC/DB) / (1 + BC/DB)
  - = GE * (CD / (DB + BC))

or Extension = w/2 * X[Bend] / (LENGTH[H][Bend] + X[Bend])

fig. 19c. Converting Wires to Boxes
The Multi-Project Chip

Insight into integrated system design is most quickly gained by actually carrying through to completion several LSI design projects, each of increasing scope. A large, complex VLSI system could be quickly and successfully developed by designers able to easily implement and test prototypes of its subsystems. The separate subsystems can be implemented, tested, debugged, and then merged together to produce the overall system layout. However, such activities are only practical if a scheme exists for carrying out implementation with minimum turnaround time and low procedural overhead per project.

In this section we describe procedures for organizing and implementing many small projects by merging their layouts onto one multi-project chip. Then each designer of a small project or subsystem need not carry the entire procedural burden involved in maskmaking and fabrication. We also include a collection of practical tips and hints that may prove useful to those undertaking their first projects or organizing their first multi project chips. While the details in this section are specific to present maskmaking and fabrication technology, they nevertheless give a feeling for the sort of things that must be done to implement projects in general. In a later section we discuss how multiple project implementation might be done in the future.

Figure 20 contains a photomicrograph of a Caltech class project chip containing 15 separate student projects. The individual projects were simply merged together onto one typically sized chip layout, approximately 3 mm by 4 mm, and implemented simultaneously as one chip type. Most of these projects are prototypes of digital subsystems designed using the methodology of this text. By implementing a small "slice" of a prototype subsystem array, one can verify that its design, layout, and implementation are correct, and measure its power and delay characteristics as yielded by the particular fabrication process, thus gaining almost as much information as would be obtained by implementing the full array.

Following fabrication, the wafers containing such multi project chips are scribed, diced, and then divided up among the participants. The typical minimum fabrication run is about 10 wafers, each ~7.5 to 10 cm in diameter. Thus even a minimum run provides a few thousand chips, and each participant ends up with many chips. Participants may then each package their chips, bonding the package leads to the contact pads of their individual project. Since most such projects are relatively small in area, yields are unusually high: if a project's design and layout have been done correctly, most of the corresponding chips will work.
Organizing a multi-project chip involves: (i) creating the layout of a *starting frame*, into which the various projects are to be merged, (ii) gathering, relocating, and merging the project layouts into the starting frame to create one design file and generating from this the PG files for the overall project chip, and (iii) documenting various parameters and specifications to be used during maskmaking and fabrication.

The starting frame contains all the auxiliary portions of the chip layout: scribe lines, alignment marks, line width testers (critical dimension marks), and test patterns. The starting frame may contain fiducial marks on each mask level if these are not to be placed by the mask house, and in some cases may contain a parity mark on each level to mark the appropriate reticle side and orientation during step and repeat reduction. A tip: placing a mask level name or symbol somewhere within the chip’s scribe line boundary on each level helps prevent the fatal error of level interchange at some time during project merging, maskmaking, or fabrication.

The contents of this starting frame must be carefully worked out to meet the requirements and constraints of the chosen mask house and fab line. The important factor of turnaround time for the entire mask and fab sequence may be reduced to some extent by repeatedly using a relatively standard starting frame which then becomes familiar to all those involved. Some typical values for the time involved: 3 to 5 weeks for maskmaking, and then 3 to 4 weeks for fabrication, longer if large work queues exist at the mask or fab firms.

When a multi-project chip is scheduled, a tentative chip partition for each project can be negotiated among the participants. Project design and layout can then proceed, with iterations on the space allocation being done right up till the final merging. The gathering and merging of project layout files into one design file is simplified if they are in a *common intermediate form*. Projects may then be relocated to their respective partitions of the chip, displayed, plotted, or otherwise checked, using minimum and consistent software operating upon manageable sized files.

When the project chip appears correctly organized, pattern generator (PG) files are produced and written on a mag tape to be sent to the mask house.

An alternative to the merging of projects at the intermediate form level, is the relocation and merging of their PG files. However, the PG files for major designs, containing fully instantiated artwork, become unwieldy in size even at today’s complexity. The PG file merging scheme is workable for projects of small to moderate size, and does provide a contingency plan for including projects having alien intermediate forms. If designs are relocated and merged at the
Fig. 20. Photomicrograph of a Caltech Class Project Chip
PG level, additional software should be provided for displaying or plotting the chip at that level, so that merging errors may be spotted. A tip: it is a good idea in any case to have some bounds checking to prevent stray items of one project from clobbering another.

A thought: the interface between design groups and mask houses would be cleaner if design files in a common intermediate form, such as CIF, rather than PG files were used to transmit designs to the patterning process. Files would be much smaller. The use of data links would be eased. The process to convert and sort design files into PG files, involving patterning mechanism dependent optimization, would be appropriately located: in association with the particular patterning mechanism.

Examples of Multi-Project Chips:

The above concepts and some further possibilities may be clarified by examining the details of some specific examples. Figure 21 illustrates a collaborative Xerox PARC/Caltech multi-project chip set [organized by D. Fairbairn, D. Johannsen, R. Lyon, J. Rowson, S. Trimmer]. The figure was produced as a software blowback from the PG file, of the metal level of this chip set. Projects in the set ranged in scope from the test of a few cells of an experimental, low power shift register [C. Sequin, U.C. Berkeley, and R. Lyon. Xerox PARC], up to a complete content addressable cache memory system [D. Fairbairn].

Although several of the projects in the set are fairly large, all were individually designed to yield chip sizes packagable in standard 40 pin packages, which can hold chips up to ~ 7 mm square. The pattern generator at the intended mask house was a GCA/D.W.Mann 3600, and the photorepeater was a Mann 3696. Together, these can produce 10x reticles having field sizes as large as 10 cm square, and can reduce, step, and repeat these at a maximum of 10mm x,y intervals onto masks. Therefore, the 3600/3696 can provide masks for square chips up to 10 mm (10,000μ) on a side. A 10mm square chip can hold the patterns of several normally sized chips. By including interior scribe lines in the starting frame, as indicated in figure 21, one reticle set can be patterned on the Mann 3600 to contain a number of different chips, each of which may contain more than one project. When masks are made, each reticle is photorepeated at intervals in x,y corresponding to its outer dimensions minus some scribe line overlap. In the example in figure 21, the x,y stepping distances were both ~9700 microns. Fabricated wafers are scribed and
diced on all scribe lines, including the interior ones, to yield chips of typical sizes. One of the projects, on the lower left chip in figure 21, is an experimental charge coupled device array [R. Davies]. The CCDs rode along on this chip set to obtain working masks for use in a completely different process technology (triple poly) from the standard nMOS the other projects used.

Figure 22 provides a higher magnification PG file software blowback of the region near the center of the left scribe line of the chip set. Alignment marks and line width testers (C/D's) were placed in this region, as noted in the figure. Software blowbacks of individual mask levels, more closely resembling the reticles and masks than would a composite design checkplot of all levels, are useful in conveying such location information to the mask and fab houses. Parity marks were not needed on the reticles for this project chip set. Fiducial marks were placed on the reticles by the mask house. Since the software converting the design files to PG files had just been constructed prior to organizing this chip set, reticle blowbacks were requested before proceeding further with maskmaking to verify that everything through pattern generation had worked correctly.

Some other practical details: Participants in the chip set shared some of the commonly used layout items normally required in any project. Examples were input contact pads with attached "lightning arrestor" circuits to protect the input MOSFET gates, and output drivers snaked around and attached to output pads. Even at current device sizes, pads occupy a large fraction of the chip area for large collections of projects, and participants tend to make the pads as small as their bonding skill allows. A square pad ~75μm on a side is a rather small bonding target, and 125μm on a side is easier for the novice to hit. Perhaps ~100μm square pads separated by ~75μm is a good compromise, and these should be at least 25μm from any other metal lines to avoid shorting the lines when bonding.

The scribe lines on this chip set were laid out as 140μm wide cuts down to 160μm wide paths on the diffusion level, to provide lanes free of oxide for scribing or sawing. Metal paths 30μm wide were then laid out straddling the boundaries of these scribe lines, to provide electrical contact from the substrate to the metal during the etching of the metal layer. Since all the projects on this chip set were prototype designs, and were not intended to be placed in extended use, the chips were not overglassed. Eliminating the overglassing meant that a mask level for defining cuts through overglassing over the contact pads and scribe lines was not needed, reducing maskmaking costs. On the other hand, the chip set included a mask level to pattern the thin gate oxide, to provide buried contacts between diffusion and poly that do not require metal coverage.
Fig. 21. Collaborative Xerox PARC/Caltech Multi-Project Chip

[Software Blowback from the PG File of the Metal Layer]
Fig. 22. Close-Up View of Figure 21 Region Containing Mask & Fab Information
Fig. 23. Multi-Project Chip Set Organized at Caltech
as does the butting contact. Such buried contacts enable more compact layouts, but are subject to a rather complex set of design rules, require an extra mask level, and sometimes reduce yield and reliability.

Deleting the overglassing process step also made it possible to electrically probe interior points on the chips during testing, probing small metal test pads included in the layouts. Such pads must be placed with care, however, because they hang relatively large capacitances onto circuitry and slow it down. Note that test pad probing requires special jigs and a stereo microscope, and that it is only possible to directly probe the metal layer. Testing uncovered chips may also require reduced light levels. The operation of dynamic circuits, i.e. those which use a pass transistor input into a gate having no other electrical connection, can be severely affected by light. Light induces leakage currents in the n-p junction between source and drain regions and the substrate. At room temperature, charge stored on dynamic nodes can be retained for many milliseconds in the absence of light. However, in normal room light the retention time is reduced to tens of microseconds. Thus care should be taken to avoid high light levels when long clocking periods are used. Dynamic memory chips are packaged in opaque black packages because of this effect.

A software blowback of the metal mask PG file of another project set, organized at Caltech, is shown in figure 23. The total area of this multi-project chip set is ~ 1 cm². It is subdivided into four major sections: The lower right quadrant contains the OM2 Data Path Chip described in Chapter 5, laid out using $\lambda = 2.5 \mu m$. The upper right quadrant contains a 16 by 16 bit multiplier with on-board accumulator [by Rod Masumoto, Caltech], also using $\lambda = 2.5 \mu m$. The lower left quadrant contains a subsystem, laid out using $\lambda = 2.9 \mu m$, which converts output from one port of a computer memory into the red, green, and blue analog signals for driving a color TV monitor. The upper left quadrant contains 28 projects, mostly from students in an LSI Systems course at Caltech. Other small projects are located along the left edge of the multiplier, and in the unused area within the TV subsystem project. The source material for this project chip set was generated on three different computer systems, in two different languages. Check plotting and viewing were done on three other systems. In addition to the Caltech projects, this chip set contains projects from Carnegie-Mellon University, Washington University (St. Louis), University of California, Irvine, and the Jet Propulsion Laboratory. Approximately 500,000 pattern generator rectangles were required to pattern the reticles for the five mask levels used in this project set. Conversion from intermediate form to PG files required ~10 CPU hours on the Caltech DECsystem 20.
The masks for the multi-project chip sets shown in figures 21 and 23 were produced by Silicon Valley mask houses from PG tapes, accompanied by PG file software blowbacks showing the locations of auxiliary layout items used during implementation, and by spec sheets containing a list of mask and fab specs and parameters. These spec sheets contain two types of information:

(i) that which the mask house will need for reading the PG tape, generating the reticles, and stepping the master masks. This includes whether dimensions are in Metric or English units, whether fiducials and parity marks have been laid out or are to be placed by the mask house, desired reticle magnification (usually 10X, sometimes 5X), the x,y step and repeat distances, the type and magnification of reticle blowbacks desired, and whether maskmaking beyond reticle generation is to be contingent upon blowback inspection. This information is independent of the chosen fab line.

(ii) that which is specific to the fab line, or lines, on which the wafers will be fabricated. Examples here are the number, size, and type of working plates desired, and the photographic polarity of the working plates, i.e. whether they are a positive or negative image of the PG pattern. The polarity of the working plates depends on the process step and on whether positive or negative resist is used. In addition, it is customary to specify how much, if any, the lines in the image will be expanded or contracted to compensate for growth or shrinkage of regions due to the process. This so-called “pulling” of line widths in maskmaking may begin as far back as at pattern generation. Thus, while the patterning and fabrication processes are design and layout independent, they are usually coupled, and masks made for a run on one fab line are not necessarily useable elsewhere.

Maskmaking and patterning technology will remain in a state of transition for years to come. The present shift is from contact printing with working plates to projection alignment using original master masks. These two alternatives are illustrated in figure 24. From the system designer’s point of view, at the interface to the mask and fab firms, they present no essential differences, requiring perhaps slightly different specs, and yielding different intermediate artifacts. In the next section we discuss the future evolution of these technologies, presenting several implementation schemes likely to become commonplace over the next decade. These schemes will enable fabrication of systems much denser and faster than present ones. However, the basic concepts of the design methodology will still apply. Remembering our film processing analogy, we will have “finer grain” and “faster” film available as time passes. However, the basic art of photography remains.
Fig. 24. Present Maskmaking and Fabrication: Two Alternatives
Patterning and Fabrication in the Future

As \( \lambda \) is scaled down toward its minimum value, ultimately limited by the physics of semiconductors to about 0.1 \( \mu \)m, it will become feasible to implement single chip, maximum density VLSI systems of enormous functional power. Patterning and fabrication at such small values of \( \lambda \) requires that certain fundamental problems be overcome\(^3\). In this section we will discuss alternative solutions to two of the major problems: At values of \( \lambda \) of \( \sim 2 \mu \)m, a problem of \textit{runout} is encountered, causing successive patterning steps to misalign over large regions of the wafers. This problem is solved by using less than full wafer exposure. At values of \( \lambda \) under 0.5 \( \mu \)m, the wavelength of light used in photolithography is too long to allow sufficient patterning resolution. This problem is solved by using non-optical lithography, exposing the resist with electron beams or x-rays.

Historically, silicon wafers have been patterned using full wafer exposure, i.e. using masks which covered the entire surface of the wafer. The pattern for one layer of one chip is stepped and repeated during the fabrication of the mask itself, so that the mask contains the patterns for a large array of chips. During the fabrication of each successive layer on the wafer, that layer's mask is aligned at two points with the pattern already on the wafer, and the entire wafer then exposed through the mask. In the future, as feature sizes are scaled down, full wafer exposure will not likely be possible for reasons developed in this section.

The earliest integrated circuits, circa 1960, were fabricated using wafers of 2.5 cm diameter, and typical chips were 1 to 2 mm, with a minimum feature size of \( \sim 25 \mu \)m. In 1978, production wafers are 7.5 to 10 cm, typical commercially manufactured LSI chips are 5 mm, and minimum feature size is \( \sim 5 \mu \). The concurrent development of ever finer features sizes and larger wafer sizes has placed an increasingly severe strain on the process of full wafer exposure. The reasons lie in the physics of wafer distortion.

When a wafer is heated to a high temperature, it expands by an amount determined by the thermal coefficient of expansion of silicon. A bare wafer will contract exactly the same amount upon cooling, and will therefore remain exactly the same size. Suppose, however, that a layer of SiO\(_2\) is grown on the wafer when it is at the high temperature. The thermal coefficient of expansion of SiO\(_2\) is approximately 1/10 that of silicon. As the wafer is cooled, the silicon will shrink at a rate much greater than that of the SiO\(_2\). Normally the resulting wafer will not be flat,
but convex on the SiO₂ side. If the wafer is cooled slowly enough, it is possible to "relieve" the stress induced by the difference in thermal contraction. Wafers in which such stress relief has been achieved are nearly flat but are, of necessity, a different size than they were originally⁶,⁷.

It might seem that subsequent masks could be scaled to just match the wafer distortion introduced up to the appropriate point in the process. Unfortunately no such correction can be introduced without a knowledge of the pattern of SiO₂ on the wafer. During cooling, dislocations are induced in the underlying silicon crystal at the edges of openings in the oxide pattern. Hence, the magnitude and direction of wafer distortion is dependent in complex ways upon the thickness and distribution of SiO₂ on the surface and upon the details of the thermal cycle. While it is in principle possible to compute a geometric correction for each pattern to be produced, it is clearly not possible to apply one correction for all possible patterns. Misalignment between subsequent layers due to distortion of this type is often referred to as runout. Runout due to wafer distortion is today the largest single contributor to misalignment between masking steps. Attempts to use finer feature sizes, which require more precise alignment, on larger wafer sizes, which induce larger distortions, seem doomed to failure unless full wafer exposure is abandoned.

Two attractive alternatives to full wafer exposure are now being explored: (i) electron beam exposure, and (ii) exposure using step and repeat of the chip pattern directly on the wafer.

A scanning electron beam system can be used to expose resist material, and is also capable of sensing a previous pattern on the surface of a wafer. The beam can initially scan an area covering the alignment marks of a particular chip. Information gained from this sensing operation can be used to compute the local distortion, and the chip can be exposed in nearly perfect alignment using these computed values. The process can be repeated for each chip on the wafer, until all have been exposed.

This technique has several virtues. No masks are required. A digital description of the chip can be exposed directly onto a silicon wafer. A different chip can be placed at each chip location, and this opens up the possibility of greatly extending the multi-project chip concept. However there are also limitations. Data is transferred serially. Even at the highest data rates which can be conveniently generated, a long time is required to expose each chip. More fundamentally, the physics of electron beam interactions places severe restrictions on the minimum practical feature size attainable. When a beam of electrons enters a resist-coated wafer, scattering occurs both in the resist and in the wafer. This backscattering contributes a partial exposure at points up to a
few microns away from the original point of beam impingement, and has a number of implications:

(i) The exposure, or spatial distribution of energy dissipation, varies with depth in the resist. Thus resist cross section is not readily controllable.

(ii) Exposure at any particular point depends on all patterns exposed within a few microns. This is known as the "cooperative exposure" or "proximity" effect and necessitates pattern-dependent exposure corrections.

(iii) Exposure latitude becomes narrower as the spatial period of a pattern is reduced. This is illustrated in figure 25, which shows the rise in background level exposure as a function of lateral distance for four different spatial periods: (a) 2 \( \mu \text{m} \), (b) 1 \( \mu \text{m} \), (c) 0.5 \( \mu \text{m} \), (d) 0.3 \( \mu \text{m} \). The beam diameter is 250 angstrom units, the energy 10keV, the resist thickness 0.4 \( \mu \text{m} \). The consequences of this background rise are particularly troublesome for high-speed, low-contrast resists. Experimental results show somewhat greater line broadening than predicted by the model.

For the above reasons, the writing time and the difficulty of exposing desired geometries increase rapidly for line-widths below about 0.5 micron.

An immediate prospect for achieving feature sizes of 1-2 \( \mu \text{m} \) with large wafers is offered by stepping the chip pattern directly on the wafer rather than on a mask. This technique avoids the serial nature of the electron beam writing by exposing an entire chip at once. Using good optical systems it has been possible for many years to produce patterns with feature sizes in the range 1 to 2 \( \mu \text{m} \). Recent progress in the design of optical projection systems may even make 1/2 to 3/4 micron line width patterns over several millimeter diameter areas practical. Techniques are known for using light to achieve alignments to a small fraction of a wavelength. Recently, an interferometric optical alignment technique has demonstrated an alignment precision of 0.02 micron and should be capable of a reregistration uncertainty less than 0.01 micron. It would seem that devices of ultimately small dimensions (0.25 \( \mu \text{m} \)) could be fabricated using optical alignment. It must be stressed that a realignment to the underlying pattern must be done at each chip location to achieve the real potential of the technique.

The step-and-align technique can be extended to ultimately small dimensions by substituting an x-ray source for the optical one, while retaining the automatic optical alignment system. X-rays require a very thin mask support, e.g. Mylar, upon which a heavy material such as gold or
tungsten is used as the opaque pattern. Interactions of x-rays with matter tend to be isolated, local events. No back-scattering of the x-rays occurs, and electrons produced when an x-ray is absorbed are sufficiently low in energy that their range is limited to a small fraction of a micron. For this reason, patterns formed by x-rays in resist materials on silicon wafers are much cleaner and better defined than those attainable by any other known technique (see figure 26). X-rays of very high intensity can be efficiently obtained from the synchrotron radiation of an electron storage ring. The time required for exposing a chip with such a source is no more than that required at present using optical exposures. Both optical and x-ray techniques have the property that the total exposure time per wafer can be made independent of how much of the wafer is exposed at a step. Therefore, the only penalty in a step and align process is the time required for mechanical motion and alignment.

It appears that we have in hand all of the techniques for ultra fine line lithography, even on larger silicon wafers. Both electron beam and optical stepping work must, however, focus on local alignment as the crucial step in achieving high density, high performance LSI.

We now describe a production lithography system for ultimately small dimensions. A major component of the system is a 500 to 700 MeV electron storage ring, approximately 5 meters in diameter, shaped in the form of a many sided polygon. The electron beam within this storage ring is deflected at each vertex by a superconducting magnet. This deflection results in a centripetal acceleration of the electrons, and hence in an intense tangential emission of synchrotron radiation. The most important component of such radiation is soft x-rays in the 280 to 1000 eV quantum energy range (wavelengths of 0.004 to 0.001 µm). Such x-rays are ideal for exposing resist materials with line widths in the 0.1µm range.12,13

One exposure station is fitted to each vertex of the storage ring. Each exposure station has an automatic optical alignment system for individual alignment of each chip.11 Coarse alignment is controlled by a laser interferometer and the wafer brought into position by ordinary lead screws moving a conventional stepping stage such as those in current photorepeaters. Auxiliary alignment features are placed on each mask level within each chip. Misalignment of two such patterns on the wafer relative to those in the mask produces Moire patterns which are detected by photosensors and fed to a computer system. Piezoelectric transducers driven by the computer system bring the wafer into final alignment under the mask. Each exposure station in such a system is capable of aligning and exposing one layer of one chip every few seconds. Each chip may contain of the order of 10^7 devices, which is the equivalent of several wafers at today's scale.
Figure 25. Electron Beam Exposure of Resist on Silicon. Monte Carlo Calculation of Exposure Level at a Silicon-Resist Interface, as a Function of Lateral Distance, for Four Spatial Periods. [contributed by H. I. Smith, Lincoln Laboratory, M.I.T.]
Figure 26. Resist on Silicon, Patterned by X-Ray Exposure.
[contributed by H. I. Smith, Lincoln Laboratory, M.I.T.]
Fig. 27 Some Future Patterning and Fabrication Alternatives
An overview of the possible routes from design files to finished chips with sub-micron layout geometries is shown in figure 27. In the immediate future, alignments much better than those achievable today will be possible with the optical step and align technique (leftmost path in figure 27). In addition, this scheme eliminates the step and repeat process in mask making, enabling considerably shorter turnaround time. The rightmost path, direct electron beam writing on the wafer, promises the ultimate in short turnaround time. It can be viewed as using the fab area as a computer output device. For high volume manufacturing, at ultimately small dimensions, the center path as described above will most likely become the workhorse of the industry.

Fully Integrated, Interactive Design Systems

{ in preparation }

- - - creating a data structure which allows the various levels of interactive processes to operate on the same data base - - - nodes, transistors, cells, and instances - - - operations on the data base - - - interactive logic transfer function tests - - - interactive circuit transfer function tests - - - interactive design rule checking - - - the filing problem - - -

System Simulation, Test Generation, and Testing

{ in preparation }

- - - system-level/register-transfer-level design description and simulation - - - testing the system design - - - practical strategies for structured VLSI system development - - - designing for testability - - - generation of test sequences - - - testing the chips - - -
References


Reading References

R1. J. J. Donovan, "Systems Programming", McGraw-Hill, 1972, an introductory text on this subject, provides practical information on the implementation and use of assemblers, macro-processors, compilers, loaders, operating systems, etc.


Chapter 5: Overview of an LSI Computer System, 
and the Design of the OM2 Data Path Chip 

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Sections: 
The OM Project at Caltech - - System Overview - - The Overall Structure of the Data Path - - The Arithmetic Logic Unit - - ALU Registers - - Buses - - Barrel Shifter - - Register Array - - Communication with the Outside World - - Data Path Operation Encoding - - Functional Specification of the OM2 Data Path Chip 

Up to this point, we have chosen simple examples to illustrate the fundamental properties of integrated systems, and have presented a design methodology which can be used to build hierarchically organized, complex systems. To more fully clarify these concepts, we now present examples drawn from the design of an LSI computer system. In this chapter, we provide a brief overview of this computer system, and then describe in detail one of its major components, the data path chip. Much of the detail in this chapter is intended to provide the reader with a source of examples of the implementation of digital logic subsystems into LSI circuit layout structures, under the constraints imposed both by the design methodology and by the architectural requirements of a real computer system. Chapter 6 similarly describes the controller chip of this computer system, and provides additional information on the sequencing of the overall system. 

In this chapter we assume that the reader is familiar with the structure and function of the classical stored program digital computer, and with the concept, and computer design implications, of microprogrammed control. An informal review of these basic concepts is given in the introductory portions of chapter 6, so that the mapping of the required controller subsystems into silicon can be examined. The less experienced reader may benefit from a study of that material in parallel with reading this chapter. 

It is important to note that the computer system discussed in chapters 5 and 6, while composed of structured LSI subsystems, is nevertheless of classical von Neumann form. The architectural possibilities of VLSI are just now beginning to be explored. Future lower cost, higher density, higher speed devices, combined with major reductions in integrated system implementation time, may make completely new forms of computing machines, and new notions of programming, not only feasible but also practical. Some of these issues will be discussed in chapter 8.
The OM Project at Caltech

The design of this computer system was undertaken as a university project in experimental computer architecture. The "Our Machine" (OM) project, as it has come to be known, was started by Carver Mead in 1976, as part of the LSI Systems course at Caltech. The project involves the design of a number of LSI chips, as described in the system overview section.

The initial focus of the project was the architecture and design of the system's primary data processing module, the data path chip. Early contributions to this effort were made by Mike Tolle [Litton Industries], while attending the LSI systems course. Other participants were Caltech students Dave Johannsen and Chris Carroll, with much inspiration from Ivan Sutherland. By December 1976, the first design (OM₀) of the data path chip was nearly completed. The participants decided at that time that the design had become "baroque" and ugly, and it was scrapped. A new data path design (OM₁) was completed by March 1977 by Dave Johannsen, Chris Carroll, and Rod Masumoto. Fabricated chips were received in June 1977. It was this chip which appeared in the September 1977 Scientific American article by Sutherland and Mead. The chip was fully functional except for a timing bug in the dynamic register array, which had been designed in departure from the structured design methodology developed in this text.

A complete redesign of the data path chip was undertaken in June 1977, by Dave Johannsen. By September 1977, a complete set of new cells had been constructed. The design was completed by December, and chips fabricated by April 1978. The redesign included improvements in the encoding of the microcode control word, and rigorously applied the structured design methodology. Certain cells from the OM₂ data path chip, and from its companion controller chip, were used as examples in chapter 3.

During 1977, the controller chip was designed as one of 4 class projects in the Caltech LSI systems course. It was finished in the summer of 1977, and fabricated chips were received in early 1978.

During 1978, the architecture of an overall system was planned. Design has begun of the three remaining chips in the OM computer system: the system bus interface chip, the memory manager chip, and the clock chip.

All of the detailed LSI design on the OM project has been done by students. Throughout most of the project's history only rather limited design aids were available, notably a simple symbolic
layout language and graphic plotters for checkplotting. The efforts of students to quickly create large integrated systems, using only primitive designs aids, helped to motivate the development and refinement of the structured design techniques described in this text.

The OM project has also required the implementation of many prototype designs and complete chip designs. Since early in the project, the Caltech group collaborated with researchers in industry, who were similarly completing many prototype LSI system designs, on the development of practical methods for simplifying and speeding up prototype project implementation. This led to the formulation and debugging of the standard starting frame for conveying multi-project chips through maskmaking and wafer fabrication, as described in chapter 4.

**System Overview**

An informal block diagram of one OM processor is shown in figure 1a. Such a processor is a complete stored program, general purpose computer. Input/output devices are usually interfaced via the external data bus and control lines, located to the left in figure 1a. Several such processors may be interconnected via the system bus to augment one user's computer system. Tasks may then be distributed among the processors, improving system performance, for example by using different processors to independently control different input/output devices. Groups of different user systems may also share the system bus.

Each OM processor is composed of five LSI chips, along with some standard memory chips and a few MSI chips. A brief description follows of the five LSI chips being designed as part of this project. For a more detailed description of these chips and the overall system, see reference 1.

The *data path chip* performs most of the data manipulation functions for the processor. These operations are performed as directed by sequences of control microinstructions, which are fetched from a microcode memory using addresses generated by the controller chip. The main subsystems of the data path chip are a register file, a barrel shifter, and an arithmetic logic unit (ALU). Two buses connect these subsystems together. This chip's internal structure is described in detail later in this chapter.

The *controller chip* contains the microprogram counter (µPC) which stores the microcode memory address, and a counter for the control of microprogram loops. This chip also contains stacks for both the microprogram counter and loop control counter values. The concepts of controller
structure and function are fundamental in computer architecture. Chapter 6 provides an introduction to these ideas, and then describes the organization and layout of the controller chip.

The *memory manager chip* provides addresses for the data memory, and directs the communication between chips on the data bus. It also implements some simple data structures in the data memory. The manager can divide the memory into separate partitions, and implement a different data structure in each partition. Four basic data structures are implemented: stacks, queues, linked lists, and arrays. When accessing a stack partition for example, the microcode need only ask the manager to push or pop data off the stack, and the manager does the rest, maintaining stack pointers, performing bounds checking to see if the stack is full or empty, etc.

The *system bus interface chip* provides asynchronous communication with other OM processors via the system bus. There are a whole host of subtleties associated with interfacing asynchronous buses. These issues are discussed in detail in chapter 7, along with the details of the organization and design of the interface chip.

The *clock chip* generates the two phase clock signals needed by the system. The clock can be stopped to allow for the synchronization of asynchronous signals. Some chips in the system have a single $q_1$ clock input, and generate the other clock phase signal on-chip.

A few words about timing may be helpful: In general, during $q_1$ data is transferred from one subsystem to another on the same chip, while during $q_2$ data is transferred from one chip to another. The data chip's ALU, and other data modification units, operate during $q_2$. Microcode is available on both phases, and is pipelined by one phase. Thus, the opcodes that control the ALU enter the data chip during $q_1$. The microprogram address is generated by the controller chip during $q_2$, gets driven off chip into the data chip's microcode latches during $q_1$, and is used to look up the next opcode on the following $q_2$. Because of these timing requirements, all jumps in the microcode are pipelined by one clock cycle.

The remainder of this chapter describes the data path chip, and is presented in two distinct parts. The first part outlines the architectural requirements for the data path chip, and then illustrates, via the detailed design and layout of the chip's subsystems and cells, how the design methodology was applied to satisfy these requirements. The second part is an external functional description of the data path chip, intended as a user manual for those who microprogram the computer system, and for reference during the study of the OM2 controller chip in chapter 6.
Figure 1a. Block Diagram of an OM Processor.

Figure 1b. General Floor Plan of the Data Path Chip.
The Overall Structure of the Data Path

The basic requirements initially established for the data path chip were (i) that it be gracefully interconnectable into multiprocessor configurations, (ii) that it effectively support a microprogrammed control structure, thus enabling machine instruction sets to be configured to the application at hand, (iii) that it be able to do variable field operations for emulation instruction decoding, assembly of bit-maps for graphics, etc., and (iv) that its performance be as fast as possible.

In order to satisfy the first requirement, the data path chip was designed with two ports: one port to be used for a system interconnection, and the other for connection to local memory, input-output devices, etc. In many systems time is lost in assembling the two operands required for many operations. Therefore, the data path has two internal buses, and all registers on the chip are two-port registers. The requirement for gracefully handling variable length words required a shifter at least sixteen bits long. The performance requirement dictated an arithmetic logic unit having considerable flexibility without sacrificing speed. In order to avoid extensive random wiring for connecting the major subsystems on the chip, the following strategy was adopted at the outset: two system buses would run through the entire processing array, from one end of the chip to the other. One port was to be located at the left end of the chip, and the other port at the right end, and the two system buses were to run the full length of the chip between the two ports through the register and the data processing array.

The three main functional blocks on the chip are the register array, the shifter, and the arithmetic logic unit. These blocks are placed next to each other in the center of the chip, between the two ports. The arrangement of the major subsystems is shown in figure 1b. The system buses run horizontally, on the polysilicon level, through these functional blocks. The major control lines run vertically across these blocks, on the metal level. The power, ground, and clock lines are run parallel to the control signal lines. The details of these functional blocks will be described in subsequent sections of this chapter. Included are descriptions of peripheral circuits needed to interface subsystems with each other and to the outside world. Detailed layouts of certain cells in the system are also included. Some of the layouts shown are earlier versions than those actually included in the final data chip. Nevertheless, they convey the basic ideas involved in laying out those cells. The overall layout of the data chip is shown in the frontispiece.
The Arithmetic Logic Unit

The carry chain of the ALU, and its associated logic, was the first functional block to be designed in detail, since it was believed that the carry chain would limit the performance of the system. Simulations of several look-ahead schemes indicated that they added a great deal of complexity to the system without much gain in performance. For this reason a decision was made early in the project to implement the fastest possible Manchester type carry chain (reference 4, chapter 1), having a carry propagation circuit similar to that shown in figure 11, chapter 1. The carry chain and its associated logic were allowed to dictate the repeat distance of the cells in the vertical direction. In MOS technology, a Manchester carry chain is particularly limited in its ability to propagate a high carry signal. However, it can quite rapidly propagate a low carry signal.

In any arithmetic logic unit there will be a null period when the OP code for the next operation is being brought in. Advantage can be taken of this null period to precharge the carry chain and other sections of the data path where timing is particularly crucial. In this way, it is not necessary to propagate high signals through pass transistors where the rise transient would be particularly slow. This strategy was applied in OM's ALU, and the resulting carry chain is shown in figure 2.

The main carry chain runs through the pass transistor from carry-in to carry-out. The carry-in signal is detected by the gate of an inverter which feeds the signal into the subsequent logic of the ALU. Three transistors are used to control the state of the carry-out of each stage. The first one merely precharges the node associated with carry-out during the null period of the ALU. The second is the carry-kill signal which is derived from the inputs to the ALU, and simply grounds the carry-out through a single transistor. The third is a pass transistor which causes carry-out to be equal to carry-in. These last two signals associated with the carry chain in each stage, carry-kill and carry-propagate, are generated by two NOR gates which have kill-bar and propagate-bar as one input and precharge as the second input. Hence, it is assured that the kill signal and propagate signal are disabled during the null period when the precharging takes place.

After some analysis, we found that nearly all interesting combinations of carry-in and the input signals could be generated using propagate and carry-in from each stage. Thus, as in fig.3, the carry-chain may be seen as a logic block with 2 inputs, carry-kill' and carry-propagate', 2 outputs, propagate and carry-in, a vertical signal, carry-in and carry-out, and one control wire, precharge.

The task of designing the balance of the ALU is now reduced to that of designing functional
Figure 2. Carry Chain Circuit for the Arithmetic Logic Unit.

Figure 3. Abstraction of the Carry Chain Circuit.
blocks to: (a) combine the two input variables to form a propagate bar and kill bar, (b) combine carry bar and propagate to form the output signal, and (c) drivers for controlling the logical function blocks and deriving a timing for precharge.

A number of random logic implementations of function blocks for deriving kill, propagate, and the output were attempted. All seemed to be at variance with the horizontally microprogrammed architecture of the data path, and required a large amount of area and power. For this reason it was decided to use the general logical function block illustrated in chapter 3, figure 12a. Recall that the depletion mode transistors, i.e. those covered by ion implanted regions represented by yellow, are always on. Such logic function blocks are used to generate carry-bar, propagate-bar, and for combining carry-bar in and propagate to form the output. The circuit, shown in figure 4, implements sixteen logic functions of two input variables. It consists of a set of transistors which fully decode the input combination of A and B, and connect one and only one of the vertical control lines to the output, depending on this input combination. For example, when A and B inputs are both low, the vertical control wire labelled $O_0$ is connected to the output. The truth table entries for the desired logic function are placed on the $G$ vertical control wires, and the output is then the desired logic function of the two input variables. For example, if the Exclusive-OR of A and B is desired, a logic-0 will be applied to the control wires 0 and 3, and logic-1 will be applied to control wires 1 and 2. Since it is desired to implement the same logic function on all bits of the word, the control variables $G_0$ through $G_3$ need not be generated in every bit slice, but may be generated once at either the top or bottom of the array. The functional abstraction of the circuit of Fig. 4 is shown in figure 5.

The block diagram for a complete arithmetic logic unit is shown in figure 6. The functional dependence of the output on the two inputs and the state of the carry is determined by a 12-bit number: $P_0$ through $P_3$, $K_0$ through $K_3$, and $R_0$ through $R_3$, together with the carry-in to the least significant bit of the ALU. The ALU is quite general, and its detailed operation set may be left unbound until the control structure of the computer system is designed at a later time.

There are two general principles illustrated by this design. First, it is often less expensive in area, time, and power to implement a general function than to implement a specific one. Secondly, if a general function can be implemented, the details of its operation can be left unbound until later, and hence, provide a much cleaner interface to the next level of design. The detailed choices of which functional entities to leave unbound and which to bind early requires a considerable amount of judgment, and is where much of the skill in integrated system design lies.
Two details need to be dealt with before the arithmetic logic unit function block is complete. Drivers are needed for the $P_0 \cdots P_3$, $K_0 \cdots K_3$, and $R_0 \cdots R_3$ control lines which will generate signals with the appropriate timing. In addition, inverters must be interposed in the carry chain occasionally to minimize the propagation delay through the entire carry chain. The way we have chosen to implement the interposition of inverters is to recognize that each carry chain function block contains two inverters which produce at their output the carry-in, having been twice inverted from the actual carry-in signal. If we merely substitute this signal for the carry-out signal from the pass transistor, we have doubly inverted our carry-in and buffered it to minimize the propagation delay. This approach avoids putting spaces between the carry function blocks for inverters. It is illustrated by the dotted connection lines in figure 2. In the actual implementation, the connection through the inverters was made in every fourth stage.

Drivers for the $P$, $K$, and $R$ control lines have the following function: At some time during the null period of the ALU (which we shall call $q_{1}$), an OP code specifying the state of each control line arrives at the drivers. It must be latched while the ALU itself is being precharged, and then it must be applied to the $P$, $K$, and $R$ control lines as soon as the ALU is activated. The $P$, $K$, and $R$ function blocks are themselves composed of pass transistors, and their outputs are more effectively driven low than high. For this reason, we will precharge the outputs of the $P$, $K$, and $R$ function blocks as well as the carry chain itself. This is most conveniently done by requiring that all of the $P$, $K$, and $R$ control signals be high during the null period of the ALU. Then, independent of the states of $A$ and $B$ inputs, the outputs will be charged high by the time the ALU active period commences. The control buffer implementing this function is shown in Fig. 7.

The OP code is latched through a pass transistor whose gate is connected to $q_{1}$, and the OP code runs into a NOR gate, the other input of which is also $q_{1}$. Thus, the output of the NOR gate is guaranteed to be low during the $q_{1}$ period. The NOR gate output is then run through an inverting super-buffer, so that during $q_{1}$ the output is guaranteed to be high. At the end of $q_{1}$, the OP code are driven onto the $P$, $K$, and $R$ control lines. The only interface specification for the ALU which must be passed to the next level of system design is that the $P$, $K$, and $R$ control signals be valid before the end of $q_{1}$, and that the $A$ and $B$ inputs likewise be valid by the end of $q_{1}$ and be stable throughout $q_{2}$, the active period of the ALU. We are then guaranteed that after enough time has passed to allow the carry to propagate, the output of the $R$ function block will accurately reflect the specified function of the ALU and may be latched at the end of $q_{2}$. 
Figure 4. General Logic Function Block Transistor Diagram.

Figure 5. Functional Abstraction of the General Logic Function Block.
Fig 4a. Stick Diagram of the Function Block

Fig 4b. Actual Layout of the Function Block

Out = G(A,B)
Figure 6. Block Diagram of a 4-Bit ALU.
ALU Registers

In order for the arithmetic logic unit described in the last section to be useful, it must be equipped with a set of registers both for its input variables and for its output. Let us consider the input registers first. Inputs to the ALU may be derived from either the shifter, the buses, or other sources. They may be latched and left unchanged during any \( q_1 \) - \( q_2 \) machine cycle or set of machine cycles. This is one of the situations in which combining the multiplexing function with the latching function simplifies the design and achieves better performance. A register operating in this manner is shown in figure 8.

The input to the first inverter can be derived from four sources: three internal sources such as shifter output, bus, etc., and a fourth, the output of the second inverter. When it is desired to latch a new signal into the register, one of the source pass transistors is driven high during \( q_1 \). The feedback transistor around the two inverters is always activated during \( q_2 \). Thus, with three vertical control wires plus the \( q_2 \) timing signal, it is possible to select one of three sources into the register, or none of the three sources, thereby leaving the previous value of the register stored on the gate of the first inverter during the \( q_1 \) period. Since it is necessary to have two inverters to form the stable pair when the feedback transistor is on, both the input and its complement are available as required by the P and K function blocks of the arithmetic logic unit. The OP code signal which selects which source will be applied to the ALU input register during \( q_1 \) must come in during the previous \( q_2 \). Each of the select signals must be low during \( q_2 \), and at most one of them may come high during the following \( q_1 \). A driver appropriate for these control signals is shown in figure 9. The control OP code is latched during \( q_2 \), during which time the NOR gate shown disables the output driver. Since the output driver in this case is non-inverting, the output select line is held low during all of \( q_2 \). At the end of \( q_2 \), the OP code signal is latched and the particular select line to be enabled that cycle is allowed to go high.

Note that this timing allows two incoming OP code bits per external wire per machine cycle. In particular, if it were desirable to share a microcode bit between the ALU function and the ALU selector inputs, this could be done by bringing the ALU OP code in during \( q_1 \) and the ALU input selection code in during \( q_2 \), as shown in figure 10. This technique was suggested by Ivan Sutherland.

The ALU output register is similar to the ALU input register, except the timing is reversed. The result of the ALU operation is available at the end of \( q_2 \).
An OP code bit will, if desired, enable the latch signal to go high during \( q_2 \). The feedback transistor is always enabled during \( q_1 \), and thus the latch is effectively static even though in the absence of a latching signal the data is stored dynamically on the gate of the first inverter through the \( q_2 \) period. Once again, both the output and its complement are available if desired.

**Buses**

An early design decision was made to have data flow through the data path chip on two buses which communicate with all of the major blocks of the system. We have already seen that the ALU performs its operation during the \( q_2 \) period and does not have valid data to place into its output register until the end of \( q_2 \). If data are to be transferred from the output register of the ALU to its input register, this must be done during the \( q_1 \) period. If we adopt a standard timing scheme in which all transfers on the buses occur during \( q_1 \), we can make use of the \( q_2 \) period when the ALU is performing its operation to precharge the buses in the same manner that the carry chain was precharged during the \( q_1 \) period. In this way we solve one of the knotty problems associated with a technology designed for ratio logic. If we had insisted that the tristate drivers associated with various sources of data for a bus be able to drive up as well as down, we would have required both a sourcing and sinking transistor, together with a method for disabling both transistors. While it is perfectly possible to build such a driver (we shall undertake the exercise as part of the design of the output ports), it is a space-consuming matter to use such a driver at every point where we wish to source data onto an internal bus. By using the bus precharge scheme, our tristate drivers become simply two series transistors as shown in figure 11.

Here the data from one source, for example the ALU output register, is placed on the gate of one of the series transistors. An enable signal which may come high during \( q_1 \) is placed on the other series transistor. If one and only one of the enable signals is allowed to come high during any one \( q_1 \) period, the bus can be driven from as many sources as necessary. The performance of such a bus is limited only by the pull-down capability of the two series transistors. We shall adopt this philosophy for the processor chip we are designing, and attach such a tristate driver to each of the output registers for the ALU.
Figure 7. ALU Control Driver
All outputs high during Phi 1 (Precharge)
Selected terms low during Phi 2
Opcode valid during Phi 1

Figure 8. ALU Input Register and Multiplexer.

Phi 1 * Select 1
Phi 1 * Select 2
Phi 1 * Select 3
Phi 2

Input 1
Input 2
Input 3

Output

Figure 9. Select Control Driver.
All outputs low during Phi 2 (Precharge)
Selected terms high during Phi 1
Opcode valid during Phi 2

Figure 10. Output Register

Phi 2 * Latch

Input

Output
Fig. 7a. ALU Control Driver Layout

Fig. 9a. Select Control Driver Layout
Figure 11. Precharged Bus Circuit.

Figure 12. A Simple 1-Bit, Right-Left Shifter.
Barrel Shifter

Since shifting is basically a simple multiplexing function, it might be thought that a shifter could be combined with the input multiplexer to the ALU. A simple 1-bit, right-left shifter implemented in this manner is shown in figure 12.

It is identical with the three-input ALU register, and the three inputs have been used to select between the bus, the bus shifted left by one, and the bus shifted right by one. To support the multibit shifts necessary for field extraction and building up odd bit arrays, something more is required. One is tempted initially to build up a multibit shift out of a number of single shifts. However, for word lengths of practical interest, the $n^2$ delay problem mentioned in Chapter 1 makes such an approach unworkable.

The basic topology of a multibit shift dictates that any bus bit be available at any output position. Therefore, data paths must run vertically at right angles to the normal bus data flow. Once this simple fact is squarely faced, a multibit shifter is seen as no more difficult than a single bit shifter. A circuit enabling any bit to be connected to any output position is shown in figure 13a. It is basically a crossbar switch with individual MOS transistors acting as the crossbar points, the basic idea being that each switch $SC_{ij}$ connects bus$_i$ to output$_j$. In principle this structure can be set to interchange bits as well as shift them, and is completely general in the way in which it can scramble output bits from any input position. In order to maintain this complete generality, the control of the crossbar switch requires $n^2$ control bits. In some applications, this $n^2$ bits may not be excessive, but for most applications a simple shift would be adequate. The gate connections necessary to perform a simple barrel shift are shown in figure 13b. The shift constant is presented on $n$ wires, one and only one of which is high during the period the shift is occurring. If the shifter's output lines are precharged in the same manner as the bus, the pass transistors forming the shift array are only required to pull down the shifter's outputs when the appropriate bus is pulled low by its tristate drivers. Thus, the delay through the entire shift network is minimized and effective use is made of the technology.

A second topological observation is that in every computing machine, it is necessary to introduce literals from the control path into the data path. However, our data path has been designed in such a way that the data bits flow horizontally while the control bits from the program store flow vertically. In order to introduce literals, some connection between the horizontal and vertical flow must occur. It is immediately obvious in figure 13b that the bus is available running vertically.
through the shift array. It is then the obvious place to introduce literals into the data path or to return values from the data path to the controller.

At the next higher level of system architecture, the shift array bit slice may be viewed as a system element with horizontal paths consisting of the bus, the shifter output, and if necessary, the shift constant since it appears at both edges of the array. The literal port is available into or out of the top edge of the bit slice, and the shift constant is available at the bottom of the bit slice. These slices, of course, are stacked to form a shift array as wide as the word of the machine being built.

One more observation concerning the multibit shifter is in order. We stated earlier that our data path was to have two buses. Therefore, in our data path, any bit slice of a shifter such as the one shown in figure 13b will of necessity have two buses running through it rather than one. We chose to show only one for the sake of simplicity. There remains the question of how the two buses are to be integrated with the shifter. Since we are constructing a two-bus data path, we have two full words available, and a good field extraction shifter would allow us to extract a word which gracefully crosses the boundary between two data path words. The arrangement shown in figure 13b performs a barrel shift on the word formed by one bus. Using the same number of control lines and pass transistors, and adding only the bus lines which are required for the balance of the data path anyway, we may construct a shifter which places the words formed by the two buses end to end and extracts a full-width word which is continuous across the word boundary between the A and B buses. This function is accomplished in as compact a form as just described with a circuit shown in figure 14. Notice that the vertical wires have a split in them. The portion of the wire above the corresponding shift output being connected to the A bus, and that below the corresponding shift output to the B bus.

It can be seen by inspection that this circuit performs the function shown in figure 15 which is just what is required for doing field extractions and variable word length manipulations. The literal port is connected directly to the A bus and may be run backwards in order to discharge the bus when a literal is brought in from the control port. A block diagram which represents the shifter at the next level of abstraction is shown in figure 16.

In order to complete the shifter functional block, it is necessary to define the drivers on the top and bottom which interface with the system at the next higher level. Let us assume that the literal bus from outside the chip will contain data which are valid on the opposite phase of the clock from that of the internal buses. In that case, a very simple interface between the two buses
Figure 13a. 4-by-4 Crossbar Switch

Figure 13b. 4-by-4 Barrel Shifter
Fig. 14. 4-by-4 Shifter with Split Vertical Wires and 2 Data Buses
Fig. 14a. Layout of a 4-Bit Barrel Shifter
A3  Shift Constant ( = 2)
A2  Shift Out 3
A1  Shift Out 2
A0  Shift Out 1
B3  Shift Out 0
B2
B1
B0

Figure 15. Conceptual Picture of the Shifter’s Operation.

Figure 16. Block Diagram of the Shifter.

Figure 17. Literal Interface.
Figure 18. A Nor Form 1-of-N Decoder.

Figure 19. A Nand Form 1-of-N Decoder.
Figure 20. A Complementary Form 1-of-N Decoder.

Figure 21. A Fully Synchronized Shifter.
which will operate in either direction is shown in figure 17.

The internal shifter output is precharged during $q_2$, and active during $q_1$. It may be sourced either from the literal bus or from the shifted combination of the A and B buses through the shift array, shown in figure 15. The external literal bus itself may be sourced either from the opposite end (the external paths from the program source) or from the end attached to the A-Bus in the shift array shown.

The bus to the external literal path is precharged during $q_1$, and data bits from the literal port of the shifter are enabled onto it by a signal active during $q_2$, as shown in Fig. 17. The two signals, $q_1 \ast \text{IN}$, and $q_2 \ast \text{OUT}$, are derived from buffers identical to those shown earlier. The shift constant itself is represented by one line out of $n$, which is high, the others remaining low. Buffers for these lines are identical to those shown in figure 9.

There is one more observation concerning the $n$-bit shift constant. It is represented most compactly by a log $n$ bit binary number. However, in order to generate from such a form a signal that can be used in the actual data path, a decoder is required to convert the binary number into a one-of-$n$ signal suitable for feeding the buffers. Decoders can be made in a number of ways in the ratio technology we are discussing. The most common form is the NOR form, which is the fully decoded equivalent of the AND-plane in the programmable logic array, Chapter 3. It is shown in figure 18. Notice that the output is a high-going one-of-$n$ pattern.

Decoders can also be made in other forms. For small values of $n$, the NAND form shown in figure 19 is often convenient. We used a variant of this form for the ALU function block described earlier. Notice that the output of this form, when used as a decoder, is a low-going one-of-$n$ pattern. There is also a complementary form of decoder which can be built with ratio technology, and was suggested by Ivan Sutherland. It takes advantage of the fact that in any decoder both the input term and its complement must be present. In this case, the input term can be used to activate pull-up transistors in series, while the complement can be used to activate pull-down transistors in parallel. This logic form is similar in principle to that used with fully complementary technologies, and has similar benefits. It can generate either a high-going or a low-going one-of-$n$ number, and dissipates no static power. A decoder of this sort is shown in figure 20. Once we have added the appropriate buffers and decoders to our shift array, we have a fully synchronized function block ready to be integrated with the system at the next level up. The properties of this block are shown in figure 21.
Register Array

In any microprogrammed processor designed for emulating an instruction set at a higher level, it is convenient to have a number of miscellaneous registers available, both for working storage during computations and for storing pointers of specific significance in the machine being emulated: stack pointers, base registers, program counters, etc. Since the data path has two buses, and the ALU is a two-operand subsystem, it is convenient if the registers in data path are two-port registers. Using the design philosophy we have been discussing, a typical two-port register cell is shown in figure 22. This register is a simple combination of the input multiplexer described earlier, the $q_2$ feedback transistor, and two tristate output drivers, one for each bus. The registers can be combined into an array $m$ bits long and $n$ bits wide, the buses passing through the array. Each cell of the array is defined at the next level up, as shown in figure 23. Drivers for the load inputs and the read outputs are identical to those shown in figure 9. While we could immediately encode the load and read inputs to the registers into $\log n$ bits, we shall delay doing so until the next level of system design. There are a number of sources for the A bus besides the registers, and we will conserve microcode bits by encoding them together.

Before we proceed, there is one mundane detail which must be taken care of in the overall topological strategy. The routing of VDD and ground must generally be done in metal, except for the very last runs within the cells themselves. Often the metal must be quite wide, since metal migration tends to shorten the life of conductors if they operate at current densities much in excess of 1 milliampere per square micron cross-section. Thus, it is important to have a strategy for routing ground and VDD to all the cells in the chip before doing the detailed layout of any of the major functional blocks. Otherwise, one is apt to be faced with topological impossibilities because certain conductors placed for other reasons interfere with the routing of the VDD and ground. A possible strategy for the overall routing of VDD and ground paths is shown in figure 24.

Notice that the VDD and ground paths form a set of interdigitated combs, so that both conductors can be run to any cell in the chip. Any strategy will do, but it must be consistent, thoroughly thought through at the beginning, and rigidly adhered to during the execution of the project.
Figure 22. A Two Port Register Cell.

Figure 23. Block Diagram Definition of the Two Port Register Cell.
Fig. 22a. Layout of Two Dual-Port Register Cells
Figure 24. VDD and GND Net for the Data Path Chip.
Figure 25. Data Port Tristate Pad Circuit

Figure 26. The Tri-State Driver, which consists of any number of Tri-State Buffer Stages followed by a Pad Driver Stage. The Current Design used Two Tri-State Buffer Stages.
Fig. 25a. Pad Driver Layout
Communication with the Outside World

Although in particular applications the interface from a port of the data path to the outside world may be a point to point communication, the ports will often connect to a bus. Thus it is desirable to use port drivers which may be set in a high impedance state. Drivers which can either drive the output high, drive the output low, or appear as a high impedance to the output are known as tristate drivers. Such drivers allow as many potential senders on the bus as necessary. Figure 25 shows the circuit for a tristate interface to a contact pad.

Here, either bus A or bus B can be latched into the input of a tristate driver during $\phi_1$. Likewise the pad may be latched into an incoming register at any time independent of the clocking of the chip. Standard tristate drivers are enabled on bus A and B. The only remaining chore is the design of the tristated buffer which drives the pad directly. Details of the tristate driver are shown in figure 26.

The terms out and outbar are fed to a series of buffer stages which provide both true and complement signals as their outputs, and are disabled by a DISABLE signal. Note that this DISABLE signal does not cause all current to cease flowing in the drivers, since the pull-up transistors are depletion type, but reduces the current to a value where it can be handled by the disable transistor of the following buffer stage. In general there will be a number of super buffer stages of this sort. The very last stage of the driver is shown in Fig. 26b. It is not a super buffer but employs enhancement mode transistors for both pull-up and pull-down. These transistors are very large in order to drive the large external capacitance associated with the wiring attached to the pad. They are disabled in the same manner as the super buffers, except that when the gates of both transistors are low, the output pad is truly tristated. Once again the two output transistors are a factor of approximately 10 larger than the last super buffer in the buffer string.

As we have seen, the inverter string necessary to transform the impedance from that of the internal circuits on chip to that sufficient for driving a pad attached to wiring in the outside world is quite large, and imposes a delay of some factor times a logarithm of this impedance ratio upon communications between the chip and the outside world. Any help which can be obtained in making this transformation is of great value. For example, the latch and buffers associated with the input bus circuit to the pad drivers can themselves be graded in impedance level, so that by the time the out and outbar signals are derived, they are at a considerably higher current drive capability than the buses. Note that the buses are a considerably larger capacitance than
minimum nodes on the chip, and thus the initial latch buffers can be larger than typical inverters on the chip. All such tricks help to minimize the number of stages between the bus and the outside pad, and thus the total delay in going off chip.

**Data Path Control Operation Encoding**

By now we have defined a complete functional data path with ports on each end and functional blocks through the center, as shown in figure 27. The data path operation code bits required to control the data path and the phase of the clock on which they are latched are shown. There are forty-nine such bits together with the four asynchronous bits for latching and driving the pad to the external world. In addition, there are the carry-out wire and the sixteen literal wires. These sixty-six wires together with the thirty-two from the left and right port must go to and come from somewhere. Schemes for encoding internal data path operations into microinstructions of various lengths are discussed in chapter 6. At one extreme all the data path control wires can be brought out to a microcode memory driven by a micro program counter and controller, in which case all operations which can be done by the data path may be done in parallel. The opposite extreme is to very tightly encode the operations of the data path into a predefined microinstruction set. In the present system, this encoding would be most conveniently done by placing a programmable logic array or set of programmable logic arrays along the top and the bottom of the data path. A condensed microinstruction could then be fed to the programmable logic arrays which would then decode the compact microinstruction into the data path operation code bits.

The important point of the design strategy we have chosen is that we can orthogonalize the design of the data path and the design of the microinstruction set in such a way that the interface between the two designs is very well defined, very clean, and can be described precisely, in a way that system designers at the next higher level can understand and work with comfortably. The data path can then be viewed as a component in the next level system design. As time progresses and it is possible to construct chips with larger and larger functional density, blocks of the sort shown will form components in even larger geometrical arrangements which will form even larger components and a whole hierarchy will emerge which will implement a system function at a much higher level than contemplated here. However, if the design strategy we have described is followed, it is possible to construct arbitrarily large and complex systems which are guaranteed to work if the individual component blocks are correct, and given the clocking period is sufficient to allow the slowest functional unit to perform its function.
Figure 27. Block Diagram of Datapath with Control Wires Added.
Using the approximate capacitance values given at the end of Chapter 2, an estimate can be made of the minimum clock period for sequencing the data path. The Phase 1 time of the data path is \( \sim 50\tau \), the same as the general estimate given in the section "Transit Times and Clock Periods" in chapter 1. However, the Phase 2 time of the data path is limited by the carry chain, as discussed earlier in this chapter. The relative areas of metal, diffusion, and gate can be estimated from the ALU layout shown in Figure 6a. The metal and diffusion occupy \( \sim 15 \) and \( \sim 8 \) times the area of the propagate pass transistor gate, respectively. Metal is \( \sim 0.1 \) and diffusion is typically 0.2 times the gate capacitance per unit area. Thus the total capacitance of each stage of the carry chain is \( \sim 4.5 \) times that of the pass transistor gate. The effective delay time is correspondingly longer than the transit time \( \tau \) of the transistor itself. The effective delay through \( n \) stages of such pass transistor logic is \( \sim \tau n^2 \). In the OM2, \( n=4 \) and the effective delay for 4 bits of carry chain is \( \sim 4.5 \times 16\tau = 72\tau \). To this must be added the delay of the doubly inverting buffers at the end of every 4 bits of straight Manchester logic. This delay is \( (1+k) \) times the transit time of the inverter pulldown, properly corrected for stray capacitance in the inverter. Here the inverter ratio \( k \) is \( \sim 8 \), since its input is driven through the pass transistors. Conservatively, strays in such a circuit are always several times greater than the basic gate capacitance, and we may estimate the inverter delays at \( \sim 30\tau \). The total carry time is thus \( \sim 100 \) times the transit time for each block of 4 ALU stages. The total Phase 2 time should then be \( \sim 400\tau \). In 1978, the fastest commercial nMOS processes yield a transit time \( \tau \) of approximately 0.3 ns, and we would expect a minimum total clock period of \( \sim 450\tau \), or \( \sim 135 \) ns.
The Second Half of this Chapter contains a functional specification of the OM2 data path chip, by Dave Johannsen of Caltech. This specification was originally documented in Display File #1111, by Dave Johannsen and Carver Mead of the Caltech Computer Science Department, and copyrighted by Caltech. The specification is reprinted here with the permission of the California Institute of Technology.
Functional Specification of the OM2 Data Path Chip

[Section contributed by David L. Johannsen, Caltech]

Introduction

This specification describes a 16-bit data path chip referred to as OM2 [986]. The OM2 contains 16 registers, an ALU, and a 32-bit shifter, and is designed as part of a micro-programmed writeable-control-store digital computer. The companion chip is the Controller chip, which contains the program counter, stacks, and so on. The Controller is described in Chapter 6. The entire system is designed to run on a single 5 volt supply.

The OM2 Datachip has two data ports for communication with the external system and a communication path to the Controller chip. The data ports are tri-state with either internal or external control. Communication with the Controller consists of a 16-bit literal port and a single flag bit. Seven control bits come directly from the microcode memory.

The system runs on a single clock, generating $\phi_1$ and $\phi_2$ internally. When the clock is high, the internal buses transfer data; when the clock is low, the ALU is performing its operation. Microcode bits enter the Datachip the phase before that code is to be executed. Therefore, the bus transfer code enters the Datachip when the clock is low, and the ALU code enters when the clock is high. Figure 1 sketches a possible OM system. For a more detailed description of system configurations, see reference 1.

Throughout this section a positive logic convention is used. A "1" refers to a high voltage level, while a "0" refers to a low voltage level.

Datapaths

A block diagram of OM2 is shown in figure 2. There are two buses which connect the various elements of the chip together. These buses transfer data while the clock is high, the period referred to as $\phi_1$. During $\phi_2$, when the clock is low, the buses are precharged. Each bus can only get data from one source, and give data to one destination during any one cycle.

The Left and Right Ports communicate between the datachip and the outside world. The
Right Port has been traditionally known as the memory bus port while the Left Port has been the system bus port, but since the two ports are identical, this is an arbitrary convention. Each port has both an input latch and an output latch to provide facilities for synchronizing the data chip to the outside buses. Under program control either of the two buses can load the output latch during \( q_1 \). There are three modes of driving data from the output latch to the pins, two of which are under program control and one of which is under hardware control. The first method is to output the data as soon as it comes from the bus, during the same \( q_1 \). The second method is to latch the data from the bus during \( q_1 \) and drive it out during the following \( q_2 \). The final method is to latch the data from the bus during \( q_1 \), but output the data when an enable pin is pulled low. The enable pin would be controlled by a bus manager, and can be asynchronous with respect to the data chip. Inputting from the port is similar. By pulling down on another enable pin, data from the external bus is loaded into the input latch, which can be read later under program control. Alternatively, the microcode can force the data currently on the external bus into the internal bus during the current \( q_1 \). With this scheme, many types of synchronous and asynchronous buses may be interfaced to OM2s. For internal control only, the external enable pins can be left floating.

Registers

The registers are static and dual port. Any one of the 16 registers may source either or both of the buses, while any one of the 16 may be the destination for either bus, but not both. There are only two restrictions to the use of the registers:

1. One register may not be the destination for both buses on the same cycle, and
2. One register may not be both the source for one bus and the destination for the other bus on the same cycle.

Shifter

The shifter concatenates the two buses, resulting in a 32-bit word, with the A bus being the more significant half. The shift constant then selects the bit position where the 16-bit output window starts. The shift constant specifies the number of bits from the B bus present in the output (i.e. a shift constant of 0 returns the A bus, while a shift constant of 15 returns the LSB of the A bus in the MSB of the output, followed by all but the LSB of
Figure 1. One Possible OM2 System Configuration
Figure 2. Block Diagram of OM2

Figure 3. Shifter Operation.
the B bus in the rest of the word). A conceptual picture of the shifter is shown in figure 3. The ALU can select as inputs either the bus, the shift output, or shift control. If shift control is selected, the entire word is 0 except where the LSB of the A bus appears in the shift output. The shifter operates on \( \eta : 1 \); it may be viewed as an extension of the buses.

**ALU**

A block diagram of a single bit of the ALU is shown in figure 4. The ALU operates on the data which is contained in its two input latches. Input latch A may be loaded from the A bus, the shifter output, or the shift control, while the input latch B may be loaded from the B bus, the shifter output, or the shift control.

The outputs of the two latches become the inputs to two function blocks which determine what will happen on the carry chain. Function block P determines whether the carry chain propagates, while K decides if it is to kill the carry. If neither are true, the carry chain generates a carry. Each function block has four control inputs, which, for the Propagate function block, are referred to as PFF, PFT, PTF, and PTT. If PFF is enabled, the P block output is high if both input latches are false (contain 0). Enabling PFT activates the output if input A is false and input B is true, and so on. If, for example, both PFF and PFT are enabled, the output is active if input A is false, regardless of the state of input B. To further illustrate the operation of the function blocks, consider addition. If both inputs contain a 1, the carry is to be generated, while if both inputs are 0, the carry is killed. If the two inputs are different, the carry is to be propagated (carry out - carry in). To do this operation, the kill output should be active if both inputs are false, so KFF is enabled. Both PFT and PTF should be enabled to propagate properly. Therefore, \( K = (KFF, KFT, KTF, KTT) = (1,0,0,0) \) and \( P = (PFF, PFT, PTF, PTT) = (0,1,1,0) \).

The result of the ALU is produced by the R function block, which has as inputs P block out and Carry in. For the addition example above, the output should be the exclusive-or of P and Cin, so \( R = (0,1,1,0) \). P, K, and R values for common ALU operations are listed in the programming section.

Two ALU output latches (A and B) can be loaded from the R block output; either one may later be used to source either bus.
Flags

The carry input to the LSB of the ALU is a logical combination of a flag bit and two control inputs. The two control inputs can force the carry in to be either 1 or 0, or they can select either flag or flag bar as the input.

There is also a method for doing conditional ALU operations under the control of a two-bit conditional OP field. A conditional operation performed by the ALU is not only a function of the control inputs, but also of the flag bit. The conditional operation control forces some of the control inputs low, regardless of what the P, K, and R microcode says. The coding for conditional operations allows the use of operations like multiply step and divide step without the necessity for branching in the microcode.

There is a 16-bit flag register which can also be a source or destination of the A bus. This register can also be loaded with the ALU flags during φ2. The ALU flags include carry out, overflow, carry in to the MSB, zero, MSB, LSB, Less than, Less than or equal to, and Higher (in unsigned value). The last three flags are comparison flags used after a subtraction. For example, after subtracting ALU input latch B from latch A, the "less than" flag is true if the value in ALU input latch B was larger than the value in ALU input latch A. The MSB of the flag register is called the flag bit, and this bit may be modified every φ1 by loading it with the value of one of the other bits of the flag register. The flag bit is used in the calculation of carry in and modification of conditional ALU Ops. This bit is also sent to the controller chip to be used for conditional branching, etc.

Literal

The one remaining datapath is the literal port. It is used to send data from the datachip to the controller, and vice versa. It is a source or destination for the A bus. When the literal port is being used, standard bus operations are suspended for that cycle.

Programming

The Datachip requires 23 bits of microcode on each phase of the clock. This section of the memo specifies the encoding of the fields within that microcode. Figure 5 shows the arrangement of the microcode word.
Figure 4. Block Diagram of one bit of the ALU

Figure 5a. Phi 2 Op Code (in on Phi 1)

Figure 5b. Phi 1 Literal Transfer Op Code (in on Phi 2)

Figure 5c. Phi 1 Normal Op Code (in on Phi 2)
Bus Transfer

The bus transfer control bits enter the datachip during \( q_2 \) and are used during the following \( q_1 \). There are two buses, the A bus and the B bus, which interconnect the modules of the Datachip. These two buses are similar in many respects; however, there are a few asymmetries as to sources and destinations. Also, when a literal is being transferred, the only bus transfer field which is active is the A bus destination, which stores the literal entered on the A bus. A listing of bus sources and destinations follows:

<table>
<thead>
<tr>
<th>A Bus Source</th>
<th>A Bus Destination</th>
</tr>
</thead>
<tbody>
<tr>
<td>00nnn Register n</td>
<td>00nnn Register n</td>
</tr>
<tr>
<td>10000 Right port pins</td>
<td>10000 Left port, drive now</td>
</tr>
<tr>
<td>10001 Right port latch</td>
<td>10001 Left port, drive ( q_2 )</td>
</tr>
<tr>
<td>10010 Left port pins</td>
<td>1001x Left port, no drive</td>
</tr>
<tr>
<td>10011 Left port latch</td>
<td>10100 Right port, drive now</td>
</tr>
<tr>
<td>10100 ALU output latch A</td>
<td>10101 Right port, drive ( q_2 )</td>
</tr>
<tr>
<td>10101 ALU output latch B</td>
<td>1011x Right port, no drive</td>
</tr>
<tr>
<td>10110 Flag register</td>
<td>11000 ALU input latch A</td>
</tr>
<tr>
<td></td>
<td>11001 ALU input latch A gets shift out</td>
</tr>
<tr>
<td></td>
<td>11010 ALU input latch A gets shift ctrl.</td>
</tr>
<tr>
<td></td>
<td>11011 Flag Register</td>
</tr>
</tbody>
</table>

<table>
<thead>
<tr>
<th>B Bus Source</th>
<th>B Bus Destination</th>
</tr>
</thead>
<tbody>
<tr>
<td>00nnn Register n</td>
<td>00nnn Left port, drive now</td>
</tr>
<tr>
<td>10000 Right port pins</td>
<td>010000 Left port, drive ( q_2 )</td>
</tr>
<tr>
<td>10001 Right port latch</td>
<td>010001 Left port, no drive</td>
</tr>
<tr>
<td>10010 Left port pins</td>
<td>01001x Right port, drive now</td>
</tr>
<tr>
<td>10011 Left port latch</td>
<td>010100 Right port, drive ( q_2 )</td>
</tr>
<tr>
<td>10100 ALU output latch A</td>
<td>010101 Right port, no drive</td>
</tr>
<tr>
<td>10101 ALU output latch B</td>
<td>01011x ALU input latch B</td>
</tr>
<tr>
<td></td>
<td>0110xx ALU input latch B gets shift output, shift const. = n</td>
</tr>
<tr>
<td></td>
<td>10nnnn ALU input latch B gets shift control, shift const. = n</td>
</tr>
</tbody>
</table>

ALU Input Selection

The two ALU input latches are destinations for the two buses, as shown in the Bus Transfer section above. In addition to being loaded directly from the buses, these two latches can be loaded from the outputs of the shift array. The shift constant always comes from the 4 least significant bits of the B Bus Destination field, even though the destination
of the B Bus is not the ALU input latch B. For example, the B Bus may be transferring the contents of register 3 into register 5 while the A Bus is transferring the contents of register 4 to the ALU input latch A through the shifter. In this case, the shift constant would be "5", because the 4 least significant bits of the B Bus Destination field contain "0101".

**ALU Operations**

The following table shows coding for ALU operations that are commonly found useful. The user is encouraged to encode other operations if these are not suitable. The numbers given are the decimal representation of the 4 bit control word. For P and K, \( A'B' = 1, AB = 2, AB' = 4, AB = 8 \). For R, \( P'C' = 1, PC = 2, PC' = 4, PC = 8 \). Cin is the carry in select, and Cond is the conditional OP select.

<table>
<thead>
<tr>
<th>K</th>
<th>P</th>
<th>R</th>
<th>Cin</th>
<th>Cond</th>
</tr>
</thead>
<tbody>
<tr>
<td>1</td>
<td>6</td>
<td>6</td>
<td>0</td>
<td>0</td>
</tr>
<tr>
<td>1</td>
<td>6</td>
<td>6</td>
<td>1</td>
<td>0</td>
</tr>
<tr>
<td>2</td>
<td>9</td>
<td>6</td>
<td>2</td>
<td>0</td>
</tr>
<tr>
<td>4</td>
<td>9</td>
<td>6</td>
<td>2</td>
<td>0</td>
</tr>
<tr>
<td>2</td>
<td>9</td>
<td>6</td>
<td>1</td>
<td>0</td>
</tr>
<tr>
<td>4</td>
<td>9</td>
<td>6</td>
<td>1</td>
<td>0</td>
</tr>
<tr>
<td>-A</td>
<td>12</td>
<td>3</td>
<td>6</td>
<td>2</td>
</tr>
<tr>
<td>-B</td>
<td>10</td>
<td>5</td>
<td>6</td>
<td>2</td>
</tr>
<tr>
<td>A + 1</td>
<td>3</td>
<td>12</td>
<td>6</td>
<td>2</td>
</tr>
<tr>
<td>B + 1</td>
<td>5</td>
<td>10</td>
<td>6</td>
<td>2</td>
</tr>
<tr>
<td>A - 1</td>
<td>12</td>
<td>3</td>
<td>9</td>
<td>2</td>
</tr>
<tr>
<td>B - 1</td>
<td>10</td>
<td>5</td>
<td>9</td>
<td>2</td>
</tr>
<tr>
<td>A &amp; B</td>
<td>0</td>
<td>8</td>
<td>12</td>
<td>0</td>
</tr>
<tr>
<td>A \lor B</td>
<td>0</td>
<td>14</td>
<td>12</td>
<td>0</td>
</tr>
<tr>
<td>A \oplus B</td>
<td>0</td>
<td>6</td>
<td>12</td>
<td>0</td>
</tr>
<tr>
<td>\neg A</td>
<td>0</td>
<td>3</td>
<td>12</td>
<td>0</td>
</tr>
<tr>
<td>\neg B</td>
<td>0</td>
<td>5</td>
<td>12</td>
<td>0</td>
</tr>
<tr>
<td>A</td>
<td>0</td>
<td>12</td>
<td>12</td>
<td>0</td>
</tr>
<tr>
<td>B</td>
<td>0</td>
<td>10</td>
<td>12</td>
<td>0</td>
</tr>
<tr>
<td>Mul</td>
<td>1</td>
<td>14</td>
<td>14</td>
<td>0</td>
</tr>
<tr>
<td>Div</td>
<td>3</td>
<td>15</td>
<td>15</td>
<td>0</td>
</tr>
<tr>
<td>A/O</td>
<td>0</td>
<td>14</td>
<td>12</td>
<td>0</td>
</tr>
<tr>
<td>Mask</td>
<td>10</td>
<td>5</td>
<td>8</td>
<td>2</td>
</tr>
</tbody>
</table>

**Carry In Select**

The Carry in select field determines what the carry into the LSB of the ALU will be,
according to the following table:

<p>| | | | |</p>
<table>
<thead>
<tr>
<th></th>
<th></th>
<th></th>
<th></th>
</tr>
</thead>
<tbody>
<tr>
<td>00</td>
<td>0</td>
<td></td>
<td></td>
</tr>
<tr>
<td>01</td>
<td>Flag bit</td>
<td></td>
<td></td>
</tr>
<tr>
<td>10</td>
<td>1</td>
<td></td>
<td></td>
</tr>
<tr>
<td>11</td>
<td>Flag bit complemented</td>
<td></td>
<td></td>
</tr>
</tbody>
</table>

**Conditional Op Select**

The conditional op select field is used to generate 3 basic conditional type operations: Multiply, Divide, and And/Or step. In a great many cases, the conditional op allows functions dependant on a flag to be performed in one cycle, rather than sending the flag to the controller and branching to two separate instructions depending upon that flag. When a conditional OP is selected, certain ALU control bits are forced to zero. Which bits are zeroed depends on the conditional OP select and the flag bit, as follows:

<table>
<thead>
<tr>
<th>Select</th>
<th>Flag bit</th>
<th>K</th>
<th>P</th>
<th>R</th>
</tr>
</thead>
<tbody>
<tr>
<td>0</td>
<td>x</td>
<td>----</td>
<td>----</td>
<td>----</td>
</tr>
<tr>
<td>1</td>
<td>0</td>
<td>----</td>
<td>----</td>
<td>----</td>
</tr>
<tr>
<td></td>
<td>1</td>
<td>----</td>
<td>----</td>
<td>----</td>
</tr>
<tr>
<td>2</td>
<td>0</td>
<td>0--0</td>
<td>0--0</td>
<td>0--0</td>
</tr>
<tr>
<td></td>
<td>1</td>
<td>0--0</td>
<td>0--0</td>
<td>0--0</td>
</tr>
<tr>
<td>3</td>
<td>0</td>
<td>----</td>
<td>----</td>
<td>----</td>
</tr>
<tr>
<td></td>
<td>1</td>
<td>----</td>
<td>----</td>
<td>----</td>
</tr>
</tbody>
</table>

**Flags**

The flag select field determines which of the ALU flags becomes the new flag bit. The following table lists the selection options.

<table>
<thead>
<tr>
<th>Select</th>
<th>New Flag Bit</th>
</tr>
</thead>
<tbody>
<tr>
<td>0</td>
<td>Old flag bit</td>
</tr>
<tr>
<td>1</td>
<td>Carry out</td>
</tr>
<tr>
<td>2</td>
<td>MSB</td>
</tr>
<tr>
<td>3</td>
<td>Zero</td>
</tr>
<tr>
<td>4</td>
<td>Less than</td>
</tr>
<tr>
<td>5</td>
<td>Less than or equal</td>
</tr>
<tr>
<td>6</td>
<td>Higher (in absolute value)</td>
</tr>
<tr>
<td>7</td>
<td>Overflow</td>
</tr>
</tbody>
</table>
The ALU flags are loaded into the flag register under the control of the latching field, bit 3.
They are loaded into the following positions:

<table>
<thead>
<tr>
<th>Bit</th>
<th>Flag</th>
</tr>
</thead>
<tbody>
<tr>
<td>0</td>
<td>Not changed</td>
</tr>
<tr>
<td>1</td>
<td>Not changed</td>
</tr>
<tr>
<td>2</td>
<td>Not changed</td>
</tr>
<tr>
<td>3</td>
<td>Not changed</td>
</tr>
<tr>
<td>4</td>
<td>Not changed</td>
</tr>
<tr>
<td>5</td>
<td>Previous value of Flag bit</td>
</tr>
<tr>
<td>6</td>
<td>Carry into MSB stage</td>
</tr>
<tr>
<td>7</td>
<td>Less than or equal</td>
</tr>
<tr>
<td>8</td>
<td>Higher (in absolute value)</td>
</tr>
<tr>
<td>9</td>
<td>Less than</td>
</tr>
<tr>
<td>10</td>
<td>LSB</td>
</tr>
<tr>
<td>11</td>
<td>Zero</td>
</tr>
<tr>
<td>12</td>
<td>MSB</td>
</tr>
<tr>
<td>13</td>
<td>Overflow</td>
</tr>
<tr>
<td>14</td>
<td>Carry out</td>
</tr>
<tr>
<td>15</td>
<td>Current Flag bit</td>
</tr>
</tbody>
</table>

**Latching Field**

The latching field specifies which of four registers should be loaded, as shown in the following table:

<table>
<thead>
<tr>
<th>Latching Field</th>
<th>Register Loaded</th>
</tr>
</thead>
<tbody>
<tr>
<td>1xxx</td>
<td>Flag register loaded with current AL flags</td>
</tr>
<tr>
<td>x1xx</td>
<td>ALU output latch A loaded with the ALU output</td>
</tr>
<tr>
<td>xx1x</td>
<td>ALU output latch B loaded with the ALU output</td>
</tr>
<tr>
<td>xxx1</td>
<td>The Literal field during the next q.2 is loaded with the contents of the A Bus during the last q.2</td>
</tr>
<tr>
<td>0000</td>
<td>None of these registers are affected</td>
</tr>
</tbody>
</table>
**Literals**

The two bit literal field specifies when a literal is to be used and which direction it goes. If both bits are 0, no literal transaction will occur. If the first bit is 1, a literal will be transferred. If the second bit is 1, the literal goes off chip, while if the bit is 0, the literal comes on chip.

**Programming Examples**

This section of the memo contains 3 programming examples which should provide a better understanding of the various datapaths within OM2.

The first example is 16-bit integer multiplication. The two inputs, X and Y, are multiplied to produce the result, Z. In the multiply loop, the number X is shifted left and the MSB is stripped off. Z is shifted left, then Y is added to the new Z if the MSB of X was a 1. The sequence of instructions is repeated 16 times, using the counter in the controller to signal when the 16 iterations have been performed. Figure 6 illustrates each step of the loop, which is listed here:

\[
\begin{align*}
\text{q2:} & \quad \text{ALU}.\text{Out}.A + \text{ALU}(\text{Shift left}) + \text{ALU}.\text{In}.A; \\
& \quad \text{Latch Flags}; \\
\text{q1:} & \quad \text{ALU}.\text{In}.A + \text{Shift out}. \quad \text{Bus}.A + \text{ALU}.\text{Out}.B; \\
& \quad \text{Bus}.B + R[1]; \quad \text{This gives a shift constant of 1.}; \\
\text{q2:} & \quad \text{ALU}.\text{Out}.B + \text{ALU}(\text{Multiply Step}); \quad \text{conditionally add}; \\
& \quad \text{Flag} + \text{Cout}; \\
\text{q1:} & \quad \text{ALU}.\text{In}.A + \text{Bus}.A + \text{ALU}.\text{Out}.A
\end{align*}
\]

The second example will be to generate a parity flag, which is not directly available from the ALU. Parity is generated by exclusive-oring all of the bits of the data together. If the data are loaded into both ALU inputs, with the B input rotated by 1, performing an exclusive-or operation will give an output that is the exclusive-or of adjacent bits; bit i of the output will be bit i of the input \( \oplus \) bit i-1 of the same input. If this same operation is performed, this time rotating the B input by 2, bit i becomes i \( \oplus \) i-1 \( \oplus \) i-2 \( \oplus \) i-3. By doing this 2 more times, rotating B first by 4 and then by 8, every bit of the output is equal to the parity: the exor of all of the bits. The MSB flag is the Parity Odd flag, while the Zero
flag is the Parity Even flag. The program is listed here, and illustrated in figure 7:

\[ \begin{align*}
\varphi_1: & \quad ALU.In.A + Bus.A + R[0]; \quad \text{generate the parity of register 0.} \\
& \quad ALU.In.B + Shift.out(1); \quad Bus.B + R[0]; \\
\varphi_2: & \quad ALU.Out.A + ALU(Exor); \\
\varphi_1: & \quad ALU.In.A + Bus.A + ALU.Out.A; \\
& \quad ALU.In.B + Shift.out(2); \quad Bus.B + ALU.Out.A; \\
\varphi_2: & \quad ALU.Out.A + ALU(Exor); \\
\varphi_1: & \quad ALU.In.A + Bus.A + ALU.Out.A; \\
& \quad ALU.In.B + Shift.out(4); \quad Bus.B + ALU.Out.A; \\
\varphi_2: & \quad ALU.Out.A + ALU(Exor); \\
\varphi_1: & \quad ALU.In.A + Bus.A + ALU.Out.A; \\
& \quad ALU.In.B + Shift.out(8); \quad Bus.B + ALU.Out.A; \\
\varphi_2: & \quad ALU(Exor); \\
\end{align*} \]

The third example adds all of the registers to what is in ALU.Out.A. By executing and modifying a literal, the registers can be indirectly accessed, which makes this routine possible. Figure 8 illustrates the operation of the following code:

\[ \begin{align*}
\varphi_1: & \quad ALU.In.A + \text{Literal} \quad "Bus.A + R[1]; \quad ALU.In.B + Bus.B + ALU.Out.B"; \\
\varphi_2: & \quad ALU.Out.B + ALU + ALU.In.A; \\
\varphi_1: & \quad ALU.In.A + Bus.A + R[0]; \\
\varphi_2: & \quad ALU.Out.B + ALU + ALU.In.A; \quad \text{This is just setup, now the loop!} \\
\varphi_1: & \quad Bus.A + ALU.Out.B; \\
& \quad ALU.In.B + Bus.B + ALU.Out.A; \\
\varphi_2: & \quad ALU.Out.A + ALU(\text{add}); \\
& \quad \text{Execute Literal;} \\
\varphi_1: & \quad ALU.In.A + A.Bus; \quad \text{The rest of this instruction is the literal!} \\
\varphi_2: & \quad ALU.Out.B + ALU(\text{increment}) + ALU.In.B; \quad \text{point to next register.}
\end{align*} \]
Figure 6a. Shift X in the ALU, putting the Cout flag into Flagbit.

Figure 6b. Put Z on Bus A, and shift 1 left in shifter.

Figure 6c. Conditionally add Z and Y.

Figure 6d. Bring X back around to the ALU input.
Figure 7a. Shifting by 1: Result is Exclusive-Or of Adjacent Bits.

Figure 7b. Shifting by 2: Result is Exclusive-Or of 4 Adjacent Bits
Figure 7c. Shifting by 4: Result is Exclusive-Or of 8 Adjacent Bits.

Figure 7d. Shifting by 8. Result Has All Bits Identically the Parity Flag.
Figure 8e. Bring Around Sum and Put Control Literal on Bus A

Figure 8f. Add Current Numbers

Figure 8g. Register Loaded by Literal Goes to ALU Input A

Figure 8h. Point to Next Register, Loop to Figure 8e


ISP Description of the OM2 Datachip

Pin States
lp < 0:17
rp < 0:17
new.code < 0:22
flag.pin < 0
power < 0:3

Pin Formats
left.port.data < 0:15
left.out.async < 0
left.in.async < 0
right.port.data < 0:15
right.out.async < 0
right.in.async < 0
literal < 0:15

clock < 0

M2 State
reg[0:15] < 0:15
a.bus < 0:15
a.bus.old < 0:15
b.bus < 0:15
left.out < 0:15
left.in < 0:15
right.out < 0:15
right.in < 0:15
left.out.later < 0
right.out.later < 0
alu.in.a < 0:15
alu.in.b < 0:15
alu.out.a < 0:15
alu.out.b < 0:15
old.code < 0:22
flags < 0:15

Instruction format
a.source < 0:4
b.source < 0:4
a.destination < 0:4
b.destination < 0:5
literal.in < 0
old.literal < 0:15
alu.p.op < 0:3
alu.k.op < 0:3
alu.r.op < 0:3
alu.conditional < 0:1
flag.select < 0:2
carry.in.select < 0:1
latch.flags < 0
latch.alu.out.a < 0
latch.alu.out.b < 0
literal.control < 0
reg.select.1 < 0:3
reg.select.2 < 0:3
reg.select.3 < 0:3

left port
right port
microcode
flag to controller
power, ground, clock, substrate

:= lp < 0:15
:= lp < 16
:= lp < 17
:= rp < 0:15
:= rp < 16
:= rp < 17
:= new.code < 5:20
:= power < 3

reg registers
bus a
bus b
left pad output latch
left pad input latch
right pad output latch
right pad input latch
for output during &2 operations
for output during &2 for right port
alu input latch a
alu input latch b
alu output latch a
alu output latch b
microcode that came in last phase
flag register

:= old.code < 5:9
:= old.code < 16:20
:= old.code < 0:4
:= old.code < 10:15
:= old.code < 22
:= old.code < 5:20
:= old.code < 19:22
:= old.code < 15:18
:= old.code < 11:14
:= old.code < 9:10
:= new.code < 6:8
:= old.code < 4:5
:= old.code < 3
:= old.code < 2
:= old.code < 1
:= old.code < 0
:= a.source < 0:3
:= a.destination < 0:3
:= b.source < 0:3
reg.select.4 < 0:3> := b.destination < 0:3>
select.1 < 0> := a.source < 4>
select.2 < 0> := a.destination < 4>
select.3 < 0> := b.source < 4>
select.4 < 0:1> := b.destination < 4:5>
shift.constant < 0:3> := b.destination < 0:3>
sharay < 0:31> := b.bus < 0:15> □ a.bus < 0:15>

Temporary State
kill.control < 0:3>
propagate.control < 0:3>
result.control < 0:3>
kill < 0:15>
propagate < 0:15>
carry < 0:16>
alu.out < 0:15>

Instruction Execution
Instruction.execution := (
  left.out.async = 0⇒(left.port.data+left.out);next
  left.in.async = 0⇒(left.in+left.port.data);next
  right.out.async = 0⇒(right.port.data+right.out);next
  right.in.async = 0⇒(right.in+right.port.data);next
  phi1 (: = clock = 1)⇒(
    left.out.later+0;next
    right.out.later+0;next
    literal.in = 1⇒(a.bus+old.literal);next
    literal.in = 0⇒(
      select.1 = 0⇒(a.bus+reg[reg.select.1]);
      select.1 = 1⇒(
        reg.select.1 = 0⇒(a.bus+right.in+right.port.data);
        reg.select.1 = 1⇒(a.bus+right.in);
        reg.select.1 = 2⇒(a.bus+left.in+left.port.data);
        reg.select.1 = 3⇒(a.bus+left.in);
        reg.select.1 = 4⇒(a.bus+alu.out.a);
        reg.select.1 = 5⇒(a.bus+alu.out.b);
        reg.select.1 = 6⇒(a.bus+flags);next);next
      select.3 = 0⇒(b.bus+reg[reg.select.3]);
      select.3 = 1⇒(
        reg.select.3 = 0⇒(b.bus+right.in+right.port.data);
        reg.select.3 = 1⇒(b.bus+right.in);
        reg.select.3 = 2⇒(b.bus+left.in+left.port.data);
        reg.select.3 = 3⇒(b.bus+left.in);
        reg.select.3 = 4⇒(b.bus+alu.out.a);
        reg.select.3 = 5⇒(b.bus+alu.out.b);next);next
      select.4 = 0⇒(reg[reg.select.4]+b.bus);
      select.4 = 1⇒(
        reg.select.4 = 0⇒(left.port.data+left.out+b.bus);
        reg.select.4 = 1⇒(
          left.out+b.bus;next
          left.out.later+1;next);
        reg.select.4 = 2⇒(left.out+b.bus);
        reg.select.4 = 3⇒(left.out+b.bus);
        reg.select.4 = 4⇒(right.port.data+right.out+b.bus));
reg.select.4 = 5 \Rightarrow 
\text{right.out} \leftarrow \text{b.bus}; \text{next} \\
\text{right.out.later} \leftarrow 1; \text{next}; \\
\text{reg.select.4} = 6 \Rightarrow (\text{right.out} \leftarrow \text{b.bus}); \\
\text{reg.select.4} = 7 \Rightarrow (\text{right.out} \leftarrow \text{b.bus}); \\
\text{reg.select.4} \in \{8,9,10,11\} \Rightarrow (\text{alu.in.b} \leftarrow \text{b.bus}); \text{next}; \\
\text{select.4} = 2 \Rightarrow (\text{alu.in.b} < 0:15> + \text{sharay} < 16\text{-shift.constant:31-shift.constant}>); \\
\text{select.4} = 3 \Rightarrow (\text{alu.in.b} + 2\text{shift.constant}); \text{next}; \text{next} \\
\text{select.2} = 0 \Rightarrow (\text{reg[reg.select.2]} \leftarrow \text{a.bus}); \\
\text{select.2} = 1 \Rightarrow 
\text{reg.select.2} = 0 \Rightarrow (\text{left.port.data} \leftarrow \text{left.out} \leftarrow \text{a.bus}); \\
\text{reg.select.2} = 1 \Rightarrow 
\text{left.out} \leftarrow \text{a.bus}; \text{next} \\
\text{left.out.later} \leftarrow 1; \text{next}; \\
\text{reg.select.2} = 2 \Rightarrow (\text{left.out} \leftarrow \text{a.bus}); \\
\text{reg.select.2} = 3 \Rightarrow (\text{left.out} \leftarrow \text{a.bus}); \\
\text{reg.select.2} = 4 \Rightarrow (\text{right.port.data} \leftarrow \text{right.out} \leftarrow \text{a.bus}); \\
\text{reg.select.2} = 5 \Rightarrow 
\text{right.out} \leftarrow \text{a.bus}; \text{next} \\
\text{right.out.later} \leftarrow 1; \text{next}; \\
\text{reg.select.2} = 6 \Rightarrow (\text{right.out} \leftarrow \text{a.bus}); \\
\text{reg.select.2} = 7 \Rightarrow (\text{right.out} \leftarrow \text{a.bus}); \\
\text{reg.select.2} = 8 \Rightarrow (\text{alu.in.a} \leftarrow \text{a.bus}); \\
\text{reg.select.2} = 9 \Rightarrow (\text{alu.in.a} < 0:15> + \text{sharay} < 16\text{-shift.constant:31-shift.constant}>); \\
\text{reg.select.2} = 10 \Rightarrow (\text{alu.in.a} + 2\text{shift.constant}); \\
\text{reg.select.2} = 11 \Rightarrow (\text{flags} \leftarrow \text{a.bus}); \text{next}; \text{next} \\
\text{flag.select.1} = 1 \Rightarrow (\text{flags} < 15> + \text{flags} < 14>); \\
\text{flag.select.1} = 2 \Rightarrow (\text{flags} < 15> + \text{flags} < 12>); \\
\text{flag.select.1} = 3 \Rightarrow (\text{flags} < 15> + \text{flags} < 11>); \\
\text{flag.select.1} = 4 \Rightarrow (\text{flags} < 15> + \text{flags} < 9>); \\
\text{flag.select.1} = 5 \Rightarrow (\text{flags} < 15> + \text{flags} < 7>); \\
\text{flag.select.1} = 6 \Rightarrow (\text{flags} < 15> + \text{flags} < 6>); \\
\text{flag.select.1} = 7 \Rightarrow (\text{flags} < 15> + \text{flags} < 13>); \text{next} \\
\text{phi2}(= \text{clock} = 0) \Rightarrow 
\text{left.out.later} = 1 \Rightarrow (\text{left.port.data} \leftarrow \text{left.out}); \text{next} \\
\text{right.out.later} = 1 \Rightarrow (\text{right.port.data} \leftarrow \text{right.out}); \text{next} \\
\text{kill.control} = \text{alu.k.op}; \text{next} \\
\text{propagate.control} = \text{alu.p.op}; \text{next} \\
\text{result.control} = \text{alu.r.op}; \text{next} \\
\text{alu.conditional} = 1 \Rightarrow 
\text{flags} < 15> = 1 \Rightarrow 
\text{propagate.control} < 0> + 0; \text{next} \\
\text{result.control} < 0> + 0; \text{next}; \text{next}; \\
\text{flags} < 15> = 0 \Rightarrow 
\text{kill.control} < 3> + 0; \text{next} \\
\text{propagate.control} < 2> + 0; \text{next} \\
\text{result.control} < 2> + 0; \text{next}; \text{next}; \\
\text{alu.conditional} = 2 \Rightarrow 
\text{flags} < 15> = 1 \Rightarrow 
\text{kill.control} < 2> + 0; \text{next} \\
\text{kill.control} < 1> + 0; \text{next} \\
\text{propagate.control} < 3> + 0; \text{next} \\
\text{propagate.control} < 0> + 0; \text{next} \\
\text{result.control} < 3> + 0; \text{next} \\
\text{result.control} < 0> + 0; \text{next};
flags < 15 > = 0 =>
   kill.control < 3 > +0;next
   kill.control < 0 > +0;next
   propagate.control < 2 > +0;next
   propagate.control < 1 > +0;next
   result.control < 2 > +0;next
   result.control < 1 > +0;next);next;
alu.conditional = 3 =>
   flags < 15 > = 1 =>
     propagate.control < 2 > +0;next
     propagate.control < 1 > +0;next);next)
kill < 0:15 > +/-
   kill.control < 3 > \(\neg\)alu.in.a < 0:15 > \(\neg\)alu.in.b < 0:15 > \(\land\)
   kill.control < 2 > \(\neg\)alu.in.a < 0:15 > \(\land\)alu.in.b < 0:15 > \(\lor\)
   kill.control < 1 > \(\land\)alu.in.a < 0:15 > \(\neg\)alu.in.b < 0:15 > \(\lor\)
   kill.control < 0 > \(\land\)alu.in.a < 0:15 > \(\land\)alu.in.b < 0:15 >);next
propagate < 0:15 > +/-
   propagate.control < 3 > \(\neg\)alu.in.a < 0:15 > \(\neg\)alu.in.b < 0:15 > \(\lor\)
   propagate.control < 2 > \(\neg\)alu.in.a < 0:15 > \(\land\)alu.in.b < 0:15 > \(\lor\)
   propagate.control < 1 > \(\land\)alu.in.a < 0:15 > \(\neg\)alu.in.b < 0:15 > \(\lor\)
   propagate.control < 0 > \(\land\)alu.in.a < 0:15 > \(\land\)alu.in.b < 0:15 >);next
carry < 0 > + carry.in.select < 1 > + (carry.in.select < 0 > \(\land\)flags < 15 > );next
for k = 1 step 1 until 16 do:
   (carry < k > + \(\neg\)(kill < k-1 > + propagate < k-1 > \(\neg\)carry < k-1 > ) + kill < k-1 > * propagate < k-1 > \(\times\) x;next
   in OM2, x is undefined)
   if kill() and propagate() are both high, the carry chain does funny things.
   We represent that here by use of the "x" in the carry function.
alu.out < 0:15 > +/-
   result.control < 3 > \(\neg\)propagate < 0:15 > \(\land\)carry < 0:15 > \(\lor\)
   result.control < 2 > \(\neg\)propagate < 0:15 > \(\land\)carry < 0:15 > \(\lor\)
   result.control < 1 > \(\land\)propagate < 0:15 > \(\neg\)carry < 0:15 > \(\lor\)
   result.control < 0 > \(\land\)propagate < 0:15 > \(\land\)carry < 0:15 >);next
latch.alu.out.a = 1 => (alu.out.a + alu.out);next
latch.alu.out.b = 1 => (alu.out.b + alu.out);next
literal.control = 1 => (literal + bus.a.old);next
latchflags = 1 =>
   flags < 5 > + flags < 15 >;next
   flags < 6 > + carry < 15 >;next
   flags < 10 > + alu.out < 0 >;next
   flags < 11 > + 0;next
alu.out = 0 => (flags < 11 > + 1);next
   flags < 12 > + alu.out < 15 >;next
   flags < 14 > + carry < 16 >;next
   flags < 13 > + flags < 14 > \(\times\) flags < 15 >;next
   flags < 9 > + flags < 12 > \(\times\) flags < 13 >;next
   flags < 7 > + flags < 11 > \(\times\) flags < 9 >;next
   flags < 8 > + \(\neg\)(flags < 14 > \(\times\) flags < 11 >);next);next);next

end of instruction execution

References:

   July 1978, Dept. of Computer Science, California Institute of Technology.

2. C. G. Bell, A. Newell, "Computer Structures: Readings and Examples", Chapter 2,
Figure 9. Pinout of the OM2 Datachip
Summary of Commands

Transfer Phase: PHI 1

<table>
<thead>
<tr>
<th>Bus A Source</th>
<th>Literal Control</th>
</tr>
</thead>
<tbody>
<tr>
<td>00000</td>
<td>Register n</td>
</tr>
<tr>
<td>10000</td>
<td>Right Port Pins</td>
</tr>
<tr>
<td>10001</td>
<td>Right Port Latch</td>
</tr>
<tr>
<td>10010</td>
<td>Left Port Pins</td>
</tr>
<tr>
<td>10011</td>
<td>Left Port Latch</td>
</tr>
<tr>
<td>10100</td>
<td>ALU Output Latch A</td>
</tr>
<tr>
<td>10101</td>
<td>ALU Output Latch B</td>
</tr>
<tr>
<td>10110</td>
<td>Flag Register</td>
</tr>
<tr>
<td>other</td>
<td>Literal (see Literal Control)</td>
</tr>
</tbody>
</table>

<table>
<thead>
<tr>
<th>Bus A Destination</th>
<th>Literal Control</th>
</tr>
</thead>
<tbody>
<tr>
<td>00000</td>
<td>Register n</td>
</tr>
<tr>
<td>10000</td>
<td>Left Port, drive now</td>
</tr>
<tr>
<td>10001</td>
<td>Left Port, drive PHI 2</td>
</tr>
<tr>
<td>1001x</td>
<td>Left Port, no drive</td>
</tr>
<tr>
<td>10100</td>
<td>Right Port, drive now</td>
</tr>
<tr>
<td>10101</td>
<td>Right Port, drive PHI 2</td>
</tr>
<tr>
<td>1011x</td>
<td>Right Port, no drive</td>
</tr>
<tr>
<td>11000</td>
<td>ALU input Latch A</td>
</tr>
<tr>
<td>11001</td>
<td>ALU input Latch A gets Shift Out</td>
</tr>
<tr>
<td>11010</td>
<td>ALU input Latch A gets Shift Control</td>
</tr>
<tr>
<td>other</td>
<td>Do Destination</td>
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</tbody>
</table>

<table>
<thead>
<tr>
<th>Bus B Source</th>
<th>Literal Control</th>
</tr>
</thead>
<tbody>
<tr>
<td>00000</td>
<td>Register n</td>
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<tr>
<td>10000</td>
<td>Right Port Pins</td>
</tr>
<tr>
<td>10001</td>
<td>Right Port Latch</td>
</tr>
<tr>
<td>10010</td>
<td>Left Port Pins</td>
</tr>
<tr>
<td>10011</td>
<td>Left Port Latch</td>
</tr>
<tr>
<td>10100</td>
<td>ALU Output Latch A</td>
</tr>
<tr>
<td>10101</td>
<td>ALU Output Latch B</td>
</tr>
<tr>
<td>other</td>
<td>No Source</td>
</tr>
</tbody>
</table>

<table>
<thead>
<tr>
<th>Bus B Destination</th>
<th>Literal Control</th>
</tr>
</thead>
<tbody>
<tr>
<td>0000000</td>
<td>Register n</td>
</tr>
<tr>
<td>0100000</td>
<td>Left Port, drive now</td>
</tr>
<tr>
<td>010001</td>
<td>Left Port, drive PHI 2</td>
</tr>
<tr>
<td>01001x</td>
<td>Left Port, no drive</td>
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<tr>
<td>010100</td>
<td>Right Port, drive now</td>
</tr>
<tr>
<td>010101</td>
<td>Right Port, drive PHI 2</td>
</tr>
<tr>
<td>01011x</td>
<td>Right Port, no drive</td>
</tr>
<tr>
<td>0110xx</td>
<td>ALU Input Latch B</td>
</tr>
<tr>
<td>0111xx</td>
<td>No Destination</td>
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<tr>
<td>1000000</td>
<td>ALU Input Latch B gets shift output, shift constant</td>
</tr>
<tr>
<td>1111111111</td>
<td>ALU Input Latch B gets shift control, shift constant</td>
</tr>
</tbody>
</table>

Operation Phase: PHI 2

<table>
<thead>
<tr>
<th>ALU Operation</th>
<th>Flag Select</th>
<th>Carry In</th>
<th>Latching Select</th>
<th>Field</th>
</tr>
</thead>
<tbody>
<tr>
<td>10000 0110 0110 00 00 Add</td>
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<td>1000 0110 0110 00 01 Add with Carry</td>
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<tr>
<td>0100 1001 0110 00 10 Subtract</td>
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</tr>
<tr>
<td>0100 1001 0110 00 01 Subtract with Borrow</td>
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<td></td>
<td></td>
</tr>
<tr>
<td>0010 1010 0110 00 00 Subtract Reversed with Borrow</td>
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<tr>
<td>0011 1100 1010 00 00 Increment A</td>
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<tr>
<td>1100 0011 0110 00 10 Increment B</td>
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<tr>
<td>0011 1010 0110 00 00 Decrement A</td>
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</tr>
<tr>
<td>0101 1010 0110 00 00 Decrement B</td>
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<tr>
<td>0000 0001 0011 00 00 Logical AND</td>
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<tr>
<td>0000 0111 0011 00 00 Logical OR</td>
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<td></td>
</tr>
<tr>
<td>0000 0110 0011 00 00 Logical Exclusive Or</td>
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<tr>
<td>0000 1100 0011 00 00 Not A</td>
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<td></td>
<td></td>
</tr>
<tr>
<td>0000 1010 0011 00 00 Not B</td>
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</tr>
<tr>
<td>0000 0111 0011 00 00 A</td>
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<tr>
<td>0000 0111 0011 00 00 B</td>
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</tr>
<tr>
<td>1000 0111 0111 01 00 Multiply Step</td>
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<tr>
<td>1100 1111 1111 10 00 Divide Step</td>
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<tr>
<td>0000 0111 0011 11 00 Conditional AND/OR</td>
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<tr>
<td>0101 1010 0001 00 10 Generate NACK</td>
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</tr>
<tr>
<td>0000 1111 10 0 0 User Defined Op</td>
<td></td>
<td></td>
<td></td>
<td></td>
</tr>
</tbody>
</table>

Carry In Select
00 o
01 Flagbit
10 1
11 Flagbit Complimented

Flag Select
000 Old Flagbit
001 Carry Out
010 MSB
011 Zero
100 Less than flag
101 Less than or equal flag
110 Higher flag
111 Overflow

Latch Field
1xxx Latch Flags
x1xx Load ALU Output Latch A
xx1x Load ALU Output Latch B
xxx Literal bits get old A Bus next PHI 1
0000 Nop

(fig10.sil)
Chapter 6: Architecture and Design of System Controllers, and the Design of the OM2 Controller Chip

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Sections:
Alternative Control Structures • • • The Stored Program Machine • • • Microprogrammed Control • • • Design of the OM2 Controller Chip • • • Examples of Controller Operation • • • Some Reflections on the Classical Stored Program Machine

This chapter presents alternative structures for controlling a data path of the type described in chapter 5. It contains a review of the basic concepts of the stored program computer, and how such computers are constructed from a combination of (i) a data processing path, (ii) a controller, and (iii) a memory to hold programs and data. A description is given of some of the ideas behind the architecture of a specific controller chip, designed at Caltech, for use with the OM2 data chip. Several examples of controller operations are provided.

We have previously used the OM2 data path chip as a source of illustrative examples, primarily at the circuit layout level, to help the reader span the range of concepts from devices, to circuit layout, to LSI subsystems. In this chapter, the controller chip is used as a source of examples one level higher, at the subsystem level, to help the reader span the range from digital logic circuits, to LSI subsystems, to arrangements of subsystems for constructing LSI computer systems. The computer system one can construct using the OM2 data chip, the OM2 controller chip, and some memory chips, contains rather simple, regular layout structures. Yet the system is functionally quite powerful, comparing well with other classical, general-purpose, stored program computers.

All present general-purpose computers are designed starting with the stored program, sequential instruction fetch-execute concepts described in this chapter. These concepts are important not only for understanding present machines, but also for understanding their limitations.

As we look into the future and anticipate the dimensional scaling of the technology, we must recognize that it will ultimately be possible to place very large numbers of simple machines on a single chip. When mapped onto silicon, classical stored program machines make heavy use of a scarce resource: communication bandwidth. They make little use of the most plentiful resource: multiple, concurrent, local processing elements. What might be the alternatives? We will reflect on some of these issues at the end of this chapter, and examine them in detail in chapter 8.
Alternative Control Structures

In this section we will clarify the distinction between the data processing and control functions in a digital computer system, and then examine several alternative forms of control structures.

The data processing path described in chapter 5 is capable of performing a rich set of operations on a stream of data supplied from its internal registers or from its input/output ports. How is it that a structure having such a static and regular appearance as the OM Data Path can mechanize such a rich set of operations? An analogy may help in visualizing the data path in operation. Imagine the data path as like a piano, with the interior regions of the chip visualized as the array of piano wires, and the control inputs along the edge of the chip as the keys. Under the external control of the controller chip, now visualized as the piano player, a sequence of keys are struck. During some cycles, many keys are struck together simultaneously, forming a chord. A complex function may thus be performed over a period of time by the data path, just as the static-looking array of piano wires may produce a complex and abstract piece of music when a series of notes and chords are struck in a particular order.

We see from this analogy, however, that the data path in itself is not a complete system. A mechanism is required to supply, during each machine cycle, the control bits which determine the function of the path during that cycle. The operations performed on data within the data path are determined by sequences of control bit patterns supplied by the system controller.

Mechanisms for supplying these sequences of control inputs to a data path can either be very simple or highly complex. There are many alternative sorts of control structures. The detailed nature of the controller has many important effects on the structure, programming, and performance of the computer system. Let us begin with the description of the simplest form of finite state machine controller. Then, through a sequence of augmentations of this controller, we will build up to the concepts of the stored program computer and microprogramming.

Simple block diagrams, such as figures 1, 2, and 3, are used in this chapter to convey the essential distinctions between various classes of controllers without requiring the diagramming of the internal details of any particular controller. Although the detailed internal logic of any particular controller may be rather complex, there are only a small set of key ideas involved in the hierarchy of controller structures presented by the sequence of block diagrams.

If you closely examine the controllers of typical computers, you will find that every one either is,
Fig.1. Finite state machine controlling the Data Path
In this case, periodically cycling through a fixed sequence of states

Fig.2. Finite state machine controlling the Data Path
In this case, the next state can be a function of the previous operation's outcome

Fig.3. Finite state machine controlling the Data Path
In this case, a data path operation result may control machine sequencing for a number of later cycles
or contains within it, a finite state machine such as those described in Chapter 3. The very simplest form of controller for the data path is a finite state machine having no inputs other than state feedback lines, as shown in figure 1. The operations performed by the data path are determined by the sequencing of the state machine. Each clock cycle, the output of the OR plane is fed back into the AND plane and determines the next state of the state machine, which periodically cycles through a fixed sequence of states. The data path is clocked in synchronism with the controller, although for simplicity we haven’t shown clock inputs to the data path in the figures. Thus a fixed algorithm implemented in the code of the state machine operates on the data in the data path.

Such a control structure could be used with the data path to implement a function such as a digital filter, in which data is taken in from the left port of the data path, a fixed set of operations performed on the data, and a result output at the right port of the data path. However, this elementary control structure provides no way to perform operations which depend on the outcome of a previous operation or upon the data itself.

A simple augmentation, shown in figure 2, enables the control sequencing to be a function of the outcome of the previous operation. In figure 2, some of the data, or some logical functions of the data, called flags, are fed into the AND plane inputs of the state machine along with the next state information. Some typical flags are: whether or not the ALU output is zero, is positive, or whether or not one ALU input is numerically equal to the other. The next state can thus be a function of flags generated during the preceding operation. To simplify figure 2, we have not shown the clock inputs to the PLA. However, assume that all subsystem structures shown in the figure, and throughout this chapter, are appropriately operated in a synchronous manner using our normal two phase clock scheme and proper design methodology.

While in principle the figure 2 structure is quite general, improvements are possible which allow greater flexibility and compactness of representation of the algorithm in the state machine. One of these improvements is shown in figure 3. Here an additional output from the OR plane of the state machine is used to control the loading of the flag outputs of the data path into a flag register. The flag register is used as an input into the AND plane of the state machine. This enables flags generated by a particular operation to be used as control inputs for the state machine for a number of later operations. The stored flag values are replaced by a new set only when the flag load signal is raised. One difficulty inherent in this structure is the limited amount of information provided by the few flags generated by the data path’s ALU.
The Stored Program Machine

A very general and powerful arrangement is shown in figure 4a. This structure is similar to the one discussed in the last section. In this case the state machine sequencing is controlled not only by the last state and flags, but also by the data coming from some memory attached to the machine. The memory contains not only the data upon which the data path is operating, but also contains encoded information for influencing the sequencing of the state machine.

This scheme gets around the limitation of the structure in figure 3, and also provides a complete new dimension of possibilities. The basic idea is to design the state machine controller so that it may perform any of a set of different predefined operations, called the machine instruction set, rather than just perform one dedicated, predefined operation. This machine instruction set is carefully defined so as to enable the system composed of the data path, controller, and memory to mechanize any of a number of different algorithms of interest to a number of different users. These algorithms are implemented as programs composed as sequences of machine instructions loaded into the memory. These programs operate upon data also contained in the memory.

It is possible to show that this arrangement is perfectly general and can implement any digital data processing function. John von Neumann\textsuperscript{1} is generally credited with originating this idea of a stored program machine, and such machines are often called von Neumann machines. The abstract notion of the most basic form of stored program machine was proposed by Turing\textsuperscript{2} in 1936, for application in the development of the theory of algorithms. The abstract Turing Machine is important not only for historical reasons, but also because of its present use in the development of the theory of computational complexity of sequential algorithms.

The way in which the stored program machine operates is as follows. One of the internal registers of the data path is selected to hold a pointer into the program stored in the memory. This register is commonly called the program counter (PC), or alternatively, the instruction address register. In one particular state of the controlling state machine, which we will call the fetch next instruction (FNI) state, the program counter is caused by the state machine to output its data as an address to the memory, and the state machine initiates a memory read from this address. The data from this memory read operation is taken into the AND plane of the state machine, placing the state machine into a state which is the first of the sequence of states which mechanize the machine instruction corresponding to the code just read from the memory. The state machine then sequences the data path through a number of specific operations sufficient to perform the
Fig. 4a. A Simple Stored Program Machine
Where data read from a memory can affect machine sequencing

Fig. 4b. A Simple Stored Program Machine
In this case, augmented by an instruction register
function defined by that instruction. At some point during instruction execution the next PC value is calculated, usually by simply incrementing the current PC value.

When the state machine has completed the interpretation, or execution, of the machine instruction, it returns to the FNI state. The instruction fetch is then repeated, sending a new program counter value to the memory as an address, reading the next instruction from the memory, and beginning its interpretation. The system can thus perform any set of required operations on data stored in memory, as specified by encoded instructions stored in memory.

There is a problem with the organization of the controller in figure 4a. Most of the steps of an instruction execution sequence need as input the encoding of the instruction which initiated the sequence. In figure 4a, this information must be duplicated each cycle by the next state information. The number of bits in the feedback path for this information can be reduced by the arrangement shown in figure 4b. Here the incoming instruction is stored in a register, called the instruction register (IR), which is loaded under the control of an output from the state machine. It stays in the instruction register, and is available for state machine input during the entire period that particular instruction is being interpreted by the machine. This new arrangement is not fundamentally different from the preceding one, but is more efficient in its use of the PLAs.

The separation and naming of the instruction register also enables us to take another step in the structuring of the state machine controller’s operations: the conception and naming of stages of the interpretation of instructions fetched and held in the IR.

Suppose we have defined a machine instruction set which, for example, includes arithmetic-logic instructions, memory instructions, and branch instructions. Suppose we also have a data path such as the OM data chip, or any other typical data path containing registers, an ALU, buses for moving data around, and inputs for control signals to control the movement of data and the ALU operations. What functions must a control unit, such as that shown in figure 4b, perform in order to fetch and execute machine instructions? We find that in most stored program machines, the execution or interpretation of each machine instruction is typically broken down into the following six basic stages. Note that some instruction types may skip one or more of the stages, and that each of the stages may require sequencing through several controller states:

1. **Fetch next instruction**: This is the starting point of the fetch-execute sequence. The machine instruction at the address contained in the PC is fetched from the memory into the IR.
(2) Decode Instruction: As a function of the fetched machine instruction's type, encoded in its OP code field, the controller must "branch" to the proper next control state to begin execution of the operations specific to that particular instruction type.

(3) Fetch instruction operands: Instructions may specify operands such as the contents of registers or of memory locations. During this execution stage, the controller cycles through a sequence of states outputting control sequences to fetch the specified operands into specified locations, for example into the input registers of the ALU.

(4) Perform Operation: The operation specified by the OP code is performed upon the operands.

(5) Store Result(s): The results of the operation are stored in destinations, such as in registers, memory locations, flags, etc.

(6) Set up next address, and return to FNI: Most instructions increment the PC by one and return to the FNI state (1). Branch instructions may modify the PC, perhaps as a function of flags, by replacing its contents with a literal value, fetched value, or computed value.

Now, how would we go about designing such a controller? We can construct the state diagram for the controller just as we did for the traffic light controller example in chapter 3. Then we proceed to build up the detailed state transition table, and finally derive the AND and OR plane code for the PLA. However, in this case the state diagram will be rather more complex than that in our earlier example. One hundred or more states may be required to implement the controller for a simple machine instruction set. How do we even begin constructing the state diagram? The above list of stages of instruction execution provides a simple means of structuring the diagram. Figure 5 contains part of the controller state diagram for a typical stored program machine. The diagram is structured as a matrix of regions, where the instruction execution stages proceed from top to bottom, and the columns contain specific state sequences for each instruction type. The FNI state is placed at the top of the diagram, followed by the states leading to the decode. The decode results in a many-way branch, each path leading to a sequence for executing a particular instruction type. The figure contains some (informal) details indicating the sorts of specific control operations performed at each stage of the instruction execution or interpretation. One will encounter many variations on the simple state diagram structure shown. These are usually easily understood elaborations. For example, groups of machine instructions may share common subsequences of control operations. To reduce the number of states, we might have another level of decoding, first decoding to groups of instructions and performing shared operations, then
Fig. 5. A Portion of the State Diagram for the Controller of a Stored Program Machine

[Illustrating some typical instruction interpretation state sequences, and their associated control outputs]
decoding to individual instruction types. In any event, the generation of the state diagram and eventually the PLA code is just a matter of grinding out the details. The generation of these details is another activity which is made more tractable by following a structured approach.

Some examples follow which will clarify how a machine instruction's execution can be divided into parts, and how the parts interact with each other. Instead of using the graphical notation of figure 5, an informal tabular form is used, containing a list of statements that are normally performed sequentially, as encountered. In these examples, the unbracketed statements under "control [& state] sequence" indicate control actions. However, some statements explicitly set the next state Y'. These statements are bracketed, "[ ]", and indicate a more complex state transition than simple state-to-state progression (shown in figure 5 by a single arrow between circles).

**ALU Example:** Suppose that an arithmetic/logic instruction in our machine instruction set has the general form: { ALUOP, REGA, REGB, REGC }, specifying that ALUOP be performed on operands REGA and REGB, and the result stored in REGC. Then the instruction { ADD, R7, R2, R5 } might be executed by the following control sequence. Note that certain of the individual control steps may occur in the same machine cycle (for example: A=R7, B=R2), as a function of the capabilities of the data path: the more the data path can do in parallel, the fewer machine cycles it will require to complete an instruction:

<table>
<thead>
<tr>
<th>Function of sub-sequence:</th>
<th>Control [&amp; State] sequence:</th>
<th>Comments:</th>
</tr>
</thead>
<tbody>
<tr>
<td>Fetch Next Inst:</td>
<td>RPORT + PC</td>
<td>Place next instr. address in right port.</td>
</tr>
<tr>
<td></td>
<td>read memory</td>
<td>Raise control line to initiate memory read.</td>
</tr>
<tr>
<td></td>
<td>PC +PC+1</td>
<td>Increment PC, overlapping incr. with fetch.</td>
</tr>
<tr>
<td></td>
<td>[Y' = fcn(memop complete)]</td>
<td>Loop here till memory read completes.</td>
</tr>
<tr>
<td></td>
<td>IR + mem data</td>
<td>Load IR with inst. when read completed.</td>
</tr>
<tr>
<td>Decode Instruction:</td>
<td>[Y' = fcn(IR)]</td>
<td>Set machine state as fcn of instruction.</td>
</tr>
<tr>
<td>Fetch Operands:</td>
<td>A = R7</td>
<td>Load ALU input registers with operands.</td>
</tr>
<tr>
<td></td>
<td>B = R2</td>
<td></td>
</tr>
<tr>
<td>Store Result:</td>
<td>R5 = ALUoutreg</td>
<td>Send result address to R5.</td>
</tr>
<tr>
<td></td>
<td>[Y' = FNI]</td>
<td>Inst. not a branch, so simply return to FNI state.</td>
</tr>
</tbody>
</table>

The example assumes there is some sort of shared access to the memory, and thus the time for completion of memory accesses is not predictable. That is why we wait, testing for the presence of a completion signal before proceeding. In some computer systems, such memory accesses might proceed in lockstep with the controller sequencing, and the data taken from, or placed on, the memory bus at some fixed number of cycles following initiation of the memory operation.
Normally, most machine instructions are not branches, so we usually just have to increment the PC sometime during instruction execution. This incrementing can often be overlapped with other control operations. In the example, the incrementing of the PC is done during the FNI stage, while waiting for the completion of the instruction-fetch memory operation.

**Memory Example:** A memory instruction in our set might have the general form \{ MEMOP, REGA, ADDRESS \}, specifying the loading or storing, according to MEMOP, of the contents of register REGA to or from the memory address ADDRESS. The instruction \{ STORE, R3, ADDRESS \} might then be executed by the following control sequence:

<table>
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<th>Control [&amp; State] sequence</th>
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</tr>
</thead>
<tbody>
<tr>
<td>Fetch Next Inst</td>
<td>RPORT + PC read memory PC + PC + 1 [Y = fcn(memop complete)] IR + mem data</td>
<td>Place next instr. address in right port. Raise control line to initiate memory read. Increment PC. Loop here till memory read completes. Load IR with inst. when read completed.</td>
</tr>
<tr>
<td>Decode Instruction</td>
<td>[Y' = fcn(IR)]</td>
<td>Set machine state as fcn of instruction.</td>
</tr>
<tr>
<td>Perform Operation</td>
<td>RPORT + IR(ADDRESS) write memory</td>
<td>Send the address of the result to the memory. Raise write control line to init. memory write.</td>
</tr>
<tr>
<td>Store Result</td>
<td>RPORT + R3 [Y' = fcn(memop complete)] [Y' = FNI]</td>
<td>Place result in right output port. Loop here till memory write completes. Inst. not a branch, so simply return to FNI state</td>
</tr>
</tbody>
</table>

**Branch Example:** Suppose that branch instructions have the form: \{ BR, COND, ADDRESS \}, specifying that if the condition COND is true according to the flags, then the PC is to be loaded from memory address ADDRESS. The branch instruction \{ BR, LT, ADDRESS \} might then be executed by the following control sequence:

<table>
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<th>Function of sub-sequence</th>
<th>Control [&amp; State] sequence</th>
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<tbody>
<tr>
<td>Fetch Next Inst</td>
<td>RPORT + PC read memory PC + PC + 1 [Y = fcn(memop complete)] IR + mem data</td>
<td>Place next instr. address in right port. Raise control line to initiate memory read. Increment PC. Loop here till memory read completes. Load IR with inst. when read completed.</td>
</tr>
<tr>
<td>Decode Instruction</td>
<td>[Y' = fcn(IR)]</td>
<td>Set machine state as fcn of instruction.</td>
</tr>
<tr>
<td>Perform Operation</td>
<td>[Y' = fcn(LT(flag))]</td>
<td>Set machine state as fcn ALU LT/flag. Set to FNI if noLT. Else continue and generate new address.</td>
</tr>
<tr>
<td>Next Address</td>
<td>PC + IR(ADDRESS) [Y' = FNI]</td>
<td>Generate new next address. Return to the FNI state.</td>
</tr>
</tbody>
</table>
Now, how are the next higher level system software control functions mapped onto this basic machine structure? Higher-level functions common to all machine instructions are often performed within the FNI stage of instruction execution. After return to the FNI state, but prior to the decode state, one machine instruction has been completely executed but no action has yet been taken to execute the next instruction. Therefore, that is a natural place to check for interrupts from I/O devices, to test the priorities for task switching in a multiprogramming environment, and so forth. The testing of these logical signals, which are input to the state machine, can often be overlapped with other FNI activity. Multiple tasks may then be implemented by having the controller manipulate multiple program-counters.

In summary, once both a machine instruction set and a data path have been defined, then the control sequences required to interpret the machine instructions can be "programmed", the overall controller state diagram constructed, the "code" for the AND and OR sections of the state machine can be generated, and software systems can be built upon the resulting stored program machine. Interestingly, the control sequences in the above examples look somewhat like "programs" written in a very primitive machine language. This observation anticipates the concept of microprogrammed control, which is described in the next section.

For more information on this material, including the various trade-offs involved in the definition and encoding of instructions, see the many examples in Bell and Newell\(^6\). See also Dietmeyer\(^7\), which works out an example all the way from state diagram through the design of the controls of an elementary digital computer. Formal methods for describing state machine algorithms are given in reference 7, and in the reference R4 of chapter 3: an interesting alternative method based on ideas of T. E. Osborne, is presented along with practical examples in Clare\(^8\).

The abstract concepts behind the arrangement shown in figure 4b are used in almost all stored program digital computers manufactured today. A computer having any sort of machine instruction set can be implemented with the arrangement shown. In many cases, the state machine is implemented in random logic and therefore is not easily recognizable as one of the forms shown. However, the operations performed are equivalent to those described here. Note that any performance contraints imposed by limited functionality in the data path simply trade off against the number of machine cycles required to mechanize particular algorithms.
Microprogrammed Control

Sometimes the complete machine instruction set is not definable at the time a computer is being designed. This contingency often arises when certain operations, defined by some later user, must be executed at very high speed. Perhaps the data path is inherently capable of satisfying the required performance constraints, but not when operated under the control of any sequence composed of standard machine instructions. In such cases, special new machine instructions would have to be defined and then implemented in the state machine control logic.

Another common situation is the need to execute the instruction set of another computer system for which the user has existing programs. While such instructions could be executed by simulation, i.e. by interpreting them via a program written in the original machine instruction set, such simulations usually pay a high performance penalty. It would be much better if the machine could execute them directly. However, a substantial augmentation and/or modification of the controller's logic would have to be made, for such direct execution to be possible.

In both of these situations it would be desirable if the state machine were implemented in some writeable medium, rather than in the fixed code of a standard programmable logic array and thus patterned permanently in the silicon. While it is quite possible to build writeable programmable logic arrays, none are currently in use. Instead, machine designers have invented many clever ways of using standard writeable memories to hold the feedback logic of the state machine.

The simplest such arrangement is shown in figure 6. Here the state machine is implemented using a set of memory chips. Collectively, this set of memory chips functions externally exactly as the programmable logic array shown earlier. However, this very elementary structure has a problem in supporting wide machine instruction words, since the decoder must exhaustively decode all combinations of the input variables. Thus, if $f$ is the number of flag bits, and $n$ is the number of next state lines, then the memory must have $2^{(i+f+n)}$ words to be of sufficient size to allow emulation of any machine having instructions $i$ bits wide. For this reason designers have taken to inserting more complex logic than just a simple instruction register into the path between the data source and the memory decoder section of state machines of this form.

A system using a logic path between the memory bus, or source of instructions, and the memory decoder section of the state machine is shown in Figure 7. Here a logic block we have termed the micro program-counter path is inserted between the source of machine instructions and the
Fig. 6. An alternative form of Stored Program Machine
Illustrating the use of a Decoder and Memory to implement
the state machine controlling the data path

Fig. 7. An alternative form of Stored Program Machine

Fig. 8. Another Way of Visualizing the Figure 7 Machine
inputs to the decoder. This type of control, using either writeable or read-only memories, is generally referred to as *microprogrammed control*. Notice in figure 7 that the flags and the machine instruction fetched from the memory both act as input data to the small micro program-counter data path, and the outputs of this data path are the microcode memory address lines. The arrangements shown is very powerful and general, and capable of emulating any instructions set for which there is sufficient microcode memory.

In a microprogrammed controller, the design of the control logic is reduced to encoding sequences of control bit patterns to be stored, along with control memory address sequencing information, in the microcode memory. The encoded control bit patterns for each clock cycle or machine cycle are visualized, as in the examples in the past section, as a primitive form of "instruction" and are called microinstructions. Rather than creating a "circles and arrows" state diagram and "assembling" PLA code, we write a symbolic microprogram and assemble it in the same manner as we would a symbolic machine language program.

The micro program-counter data path ($\mu$PC) is similar to the main data path: it is controlled by a number of outputs from the microcode memory section of the state machine. Its main purpose is to decrease the amount of microcode memory required to emulate the particular machine instruction set being implemented. This is done in two ways: First, the $\mu$PC maps the $f+n$ bits of state into a smaller number of bits which are then decoded to address the microcode memory. Secondly, it reduces $n$ by allowing complex operations within the $\mu$PC to be specified with only a few bits of control information. The controller chip described in the later sections of this chapter is the microprogram counter path portion of a microprogrammed controller for OM2.

The concept of microprogramming was originated by M. V. Wilkes$^{3,4}$ in 1951. In those days when controller logic functions were implemented using gates constructed out of vacuum tubes, switching hardware was very expensive compared to wires, and great efforts were expended towards gate minimization. This inevitably led to rather intertwined connections in the controller logic, and any change in function might require a complete redesign. Wilkes presented the notion of microprogrammed control using a read-only memory to hold the control sequences, as a means of bringing regularity and structure to the design of system controllers and thus simplifying their design and redesign. There is a large body of knowledge associated with the architectural implications of microprogrammed control, and the serious reader will benefit from a study of the literature$^{5,6,7}$. 
Today, although we can easily implement control logic in a structured way using a PLA, we still often use microprogrammed control in order to obtain the advantages offered by writeable control logic. An additional present advantage of microprogrammed control is that the detailed design/redesign of control logic is extended into the wide arena of those familiar with linear sequential programming concepts. In the future as the "programming" of structures into silicon becomes easier, as the time to implement designs becomes much shorter, and as state machine "coding" becomes more widely understood, we may find that these activities will become viewed as a natural extension of microprogramming.

There is an alternative way of viewing the machine shown in figure 7. Examine carefully the loop formed by the micro program-counter data path, the decoder section of the microcode memory, and the outputs of the microcode memory which are used to control the micro program-counter. We can view the microcode memory address as an instruction address and the wires coming from the microcode memory to control the micro program-counter path as an instruction.

This alternative view is illustrated in figure 8. Observe that we have constructed another stored programmed machine of the same form as that shown in figure 4b. We have come full circle in our machine design: in our zeal to put as much capability as possible in the path between the machine instruction and the decoder of the state machine, we have in fact created a stored programmed machine within a stored programmed machine. This phenomenon is referred to by Ivan Sutherland as the "great wheel of reincarnation". Computers often have many such levels of machine within them, each a general purpose stored program machine in its own right. We thus find that elaborate computing machines are often only simple machines, nested and connected together in complex ways.
Design of the OM2 Controller Chip

We now describe some of the ideas behind the design of one particular micro program-counter path used for controlling the OM data chip in the system configuration\(^9\) described in chapter 5. The design of the controller chip will be examined at several stages in its actual development. This material illustrates the mapping into LSI, and the topological/geometrical planning in LSI of various subsystems such as stacks, incrementers/decrementers, multiplexers, etc., which are useful in constructing controllers.

Even at the 1978 value of \(\lambda = 3\) microns, the OM2 data path and certain forms of controller could be integrated onto a single chip. The separation of these modules onto two chips was primarily for research and tutorial purposes in the university environment: so that different controllers could be used with the OM2 data chip and vice-versa. The fact that data path and controller are on separate chips does, however, lead to detailed system partitioning decisions aimed at minimizing interchip communication. These decisions might be made differently were data path and controller integrated onto the same chip. Nevertheless the issue of minimization of interchip communication would still be involved at the next system level, and is worthy of study.

The basic function of the *micro program-counter path*, which we call the *controller* for short, is to provide microprogram memory addresses. The microprogram memory addresses are stored in a latch which is called the *micro program-counter*, or \(\mu\)PC. The \(\mu\)PC should be distinguished from the *program-counter*, or PC, which stores the main memory addresses of higher level machine instructions. The most common address calculation is to increment the address by one, so in addition to the \(\mu\)PC latch, the controller should contain an incrementer. The second most important address calculation is the jump or branch, so there should be some means of forcing values into the \(\mu\)PC latch. With the hardware mentioned so far, we have progressed one step beyond the controller type shown in figure 6; our instruction register also increments, so we don’t need the feedback terms that originate in the microcode memory and drive the memory decoder.

A great deal of microcode memory space can be saved if *subroutines* are available at the microcode level. These subroutines can be shared between microcode sequences emulating instructions at a higher level. For example, many different machine instruction types may have the same set of operand fetch sequences. If the machine instruction set encodes a variety of indexing or relative addressing schemes, these operand fetch sequences may be quite lengthy, and repeating these sequences for every instruction type would waste a great deal of microcode
memory. To provide such microcode subroutine capabilities, provisions must be made for saving \( \mu \text{PC} \) values, which is most easily done with a stack. Stacks are easily constructed in LSI. An example of stack cell and subsystem design, and stack control driver design, is given in chapter 3.

The microcoder may also wish to use relative jumps or subroutine calls so that relocatable microcode can be written. To provide for relative operations, an adder must be included that can add displacements to the \( \mu \text{PC} \) contents. The displacements can either be fixed displacements and come from the microcode or be calculated displacements and come from the data path. Calculated displacements enable many-way branching, or dispatching, in the microcode, which is an almost essential operation for emulating instructions at a higher level. An example of dispatching will be given in a later section. Therefore, provisions should be made for accepting displacements from either the microcode or the data path.

Another microcode address operation that could be considered is a form of loop operation, which is useful when sections of microcode should be executed \( n \) times, where \( n \) can either be a constant and come from the microcode or be the result of a calculation done in the data path. One way to implement this instruction is to dedicate one register in the data path to be the loop counter and to do conditional branches in the controller based on the result of decrementing the value in that register. This is simple to do, because the hardware of the controller and data path discussed so far will allow the execution of this instruction. Unfortunately, there is a time penalty when doing interchip communication: the loop counter must be decremented during one cycle, the result of the decrement must be sent to the controller during the following cycle, and a conditional branch must be performed in the controller on the third cycle. If the loop counter were in the controller chip, this operation would only take one cycle and would not require the use of the ALU for one cycle in the data path.

With only one loop counter, loops could not be nested, and loops could not be used inside of subroutines. If a stack were provided for the counter values, however, nested loops and loops within subroutines could both be accommodated.

The first OM2 controller proposal was based primarily on the arguments presented above. Figure 9 shows a block diagram of the proposed controller. Table 1 lists the operations possible for each of the three sections of the controller chip: the \( \mu \text{PC} \) source selection, the \( \mu \text{PC} \) stack operation, and the loop stack operation. In each cycle, the controller executes one operation in each of the three sections. For most operations, all three sections work together to perform the programmed
Figure 9. Initial Block Diagram of the Controller Chip.

Figure 10. Final Block Diagram of the Controller Chip.
### Table 1. Opcodes of the Initial Controller Proposal.

<table>
<thead>
<tr>
<th>uPC Sources:</th>
<th>uPC Stack Operation:</th>
</tr>
</thead>
<tbody>
<tr>
<td>uPC + 1</td>
<td>Push uPC + 1</td>
</tr>
<tr>
<td>microcode</td>
<td>Push microcode</td>
</tr>
<tr>
<td>uPC Stack Top</td>
<td>Push uPC + microcode</td>
</tr>
<tr>
<td>True: uPC Stack Top; False: uPC + 1</td>
<td>Push uPC Stack Top + microcode</td>
</tr>
<tr>
<td></td>
<td>Pop</td>
</tr>
<tr>
<td></td>
<td>Push uPC + literal</td>
</tr>
<tr>
<td></td>
<td>True: Pop; False: NOP</td>
</tr>
<tr>
<td></td>
<td>True: NOP; False: Pop</td>
</tr>
<tr>
<td>Loop Stack Operation:</td>
<td>NOP</td>
</tr>
<tr>
<td>Push microcode</td>
<td></td>
</tr>
<tr>
<td>Push literal</td>
<td></td>
</tr>
<tr>
<td>Push count</td>
<td></td>
</tr>
<tr>
<td>Decrement Count</td>
<td></td>
</tr>
<tr>
<td>NOP</td>
<td></td>
</tr>
</tbody>
</table>

### Table 2. Opcodes of the Final Controller Design.

<table>
<thead>
<tr>
<th>uPC Sources:</th>
<th>uPC Stack Sources:</th>
</tr>
</thead>
<tbody>
<tr>
<td>uPC + 1</td>
<td>Adder output</td>
</tr>
<tr>
<td>uPC + microcode + 1</td>
<td>uPC</td>
</tr>
<tr>
<td>microcode</td>
<td>microcode</td>
</tr>
<tr>
<td>Stacktop + 1</td>
<td>literal</td>
</tr>
<tr>
<td>Stacktop + microcode + 1</td>
<td></td>
</tr>
<tr>
<td>Stacktop + literal + 1</td>
<td></td>
</tr>
<tr>
<td>uPC + literal + 1</td>
<td></td>
</tr>
<tr>
<td>literal + microcode</td>
<td></td>
</tr>
<tr>
<td>Condition Selection:</td>
<td>Counter Operations:</td>
</tr>
<tr>
<td>False</td>
<td>No Operation</td>
</tr>
<tr>
<td>True</td>
<td>Push microcode</td>
</tr>
<tr>
<td>Data path flag</td>
<td>Push literal</td>
</tr>
<tr>
<td>Compliment of Data path flag</td>
<td>Pop to literal bus</td>
</tr>
<tr>
<td>Count = 0</td>
<td>true: decrement; false: pop</td>
</tr>
<tr>
<td>Count &lt; 0</td>
<td>true: decrement; false: nop</td>
</tr>
<tr>
<td>Data path flag AND Count = 0</td>
<td></td>
</tr>
<tr>
<td>Data path flag OR Count = 0</td>
<td></td>
</tr>
<tr>
<td>uPC Stack Operations:</td>
<td>Push if condition is true</td>
</tr>
<tr>
<td></td>
<td>No push</td>
</tr>
<tr>
<td></td>
<td>Pop if condition is true</td>
</tr>
<tr>
<td></td>
<td>No pop</td>
</tr>
</tbody>
</table>
operation. There are cases, however, when only one or two of the sections are needed to perform the controller's instruction, so the other section(s) are free to perform other tasks. For example, the loop stack may be loading a count from the Data Path, while the \( \mu \text{PC} \) sections are performing a subroutine call. This concurrency saves having to load the count later, and may save microcode space. Because the controller's instruction is broken into three fields, more than one thing can be happening in parallel in the controller. This is why the instruction was not kept as one field and decoded into the three sections on chip.

The controller shown could handle all of the microcode address operations listed above, and a few new operations were discovered and added to the list. However, there are a few problems with this design. It is a "brute force" design: rather than viewing the whole chip at one time and looking for generalizations, each section of the chip and of the chip's operation was looked at individually and the chip was filled with specialized hardware for performing specialized operations. It was found that by adding one circuit here a new operation could be performed, and that by adding another circuit there a different operation could be added to the repertoire. Many designs suffer from "creeping features" of this sort. While it may be easy to draw circles and arrows on paper, it can be more difficult to draw adders and multiplexers on silicon. It would be very difficult to route all the wires needed to interconnect the devices shown in the proposal.

So let's make a few generalizations about the circuits in the design. First, there are too many adders on the chip. A close look at the proposal shows that for almost all operations we only use one adder for any one cycle, and the few operations that used more than one adder are not critical operations. Incrementing the \( \mu \text{PC} \) can also be done in the adder, by clearing one of the data inputs to the adder and forcing a carry into the first stage. Thus, all three of the adders and the increneter can be combined into one adder, and multiplexers can be put on the inputs to that adder. Another simplification would be to always load the \( \mu \text{PC} \) latch from the output of the adder, which would allow the removal of the multiplexer on the input to the latch. The only operations that were sacrificed in making the simplifications involved loading the \( \mu \text{PC} \) stack with the output of an adder. Figure 10 shows the block diagram of our simplified controller, and Table 2 lists the operations it performs. Notice that the controller's instruction is now broken into five fields, controlling the \( \mu \text{PC} \) sources, the \( \mu \text{PC} \) stack sources, the counter operation, the condition selection, and the \( \mu \text{PC} \) stack operation.

Now we will develop the geometrical and topological arrangement of the controller's subsystems.
Such arrangements are often called floor plans. A translation of the preceding ideas into the starting floor plan of the controller is shown in figure 11. The plan is composed of subsystems built of horizontal bit slices which are then stacked vertically. The number of bit slices is equal to the microcode address width for the machine, which in this case is 12 bits.

The following points were considered when deciding upon the basic framework of the floor plan. First, the \( \mu \) PC latch is placed adjacent to the microcode memory address pins. This is done to minimize the delay when driving addresses to the memory, as this operation is in the critical timing path for the entire machine. The input of the latch comes only from the output of the adder, so the adder should logically be placed next to the \( \mu \) PC latch. The adder is considerably simpler than the full arithmetic logic unit used in the data path. However, it employs the same principles as the ALU: the Manchester carry chain, the insertion of double inverters every four bits to minimize the delay in the carry chain, and the logic block to implement the desired functions with the minimum delay and power. The multiplexer is placed adjacent to the left side of the adder. This multiplexer operates in the same manner as the input multiplexer to the ALU in chapter 5. The \( \mu \) PC stack is then placed to the left of the multiplexer.

The only problem with this arrangement of the floor plan is that the microcode bus and the data path bus must also connect to the multiplexer. A large amount of area would be wasted if these two buses connected to the multiplexer from the side. Instead, if the buses could be placed where the \( \mu \) PC stack is located, they could connect to the loop counter circuits directly. But then there is the problem of where to place the \( \mu \) PC stack! One solution is to run the buses through the \( \mu \) PC stack! Each cell of the stack thus has the two buses designed right in. The two buses could then run on through the loop counter stack to the loop counter decremented and the pads.

Having placed the major blocks of the chip into the floor plan, the layout of the control circuits can be examined, and a detailed floor plan worked out. Each of the stacks require push and pop drivers, as discussed in Chapter 3. As in the chapter 3 example, one set of drivers is placed along the top, and the other set along the bottom of the stack. The control drivers for the latch, adder, multiplexer, and counter are identical to those discussed in Chapter 5. The control bits for these control drivers could all be derived directly from the outputs of the microcode memory, but this technique would result in an exceedingly wide microinstruction. By encoding the operations to be performed by the adder and its input multiplexers, the width of the microinstruction can be dramatically decreased. With proper encoding of these operations, the functional capability of the chip is not impaired, since a number of possible control signal combinations are in fact illegal and
Fig. 11. Starting Floor Plan for the Controller

Fig. 12. Final Floor Plan of the Controller Chip
thus redundant. For instance, if more than one control line for the multiplexer is enabled, the outputs of two or more sources would be shorted together, and the resulting multiplexer output would contain erroneous data. The placements of the control circuits and encoding PLAs are shown in figure 12, which also shows additional details of the final floor plan. Notice that the counter stack is higher than the 12-bit high μPC stack, so that it can contain entire 16 bit data path words for parameters passed to subroutines in the microcode. The stacks are aligned on their least significant bit position, and the additional length of the counter stack allowed space for the control PLAs for the adder and μPC stack.

The programmable logic arrays employed in instruction decoding do not have feedback from their outputs back into their inputs. Their only function is to serve as combinational logic for condensing the number of control wires and thus saving microcode memory bits. The finite state machine for the control of this path is made up of the microcode memory address feedback through the adder and stack PLAs and also the microcode literal path feedback into the input of the adder. If there were feedback terms in any of the PLAs, provisions would have to be made for access to the state of the feedback terms from off chip. Without such access, the untestable state information on the chip would make the testing of the completed chip next to impossible: the current operation of the chip would be a function not only of the control signals and data that we supply to the chip at a particular moment, but also of the past control signals and data. In the absence of a practical way to directly probe all the signal lines on the chip, it is imperative that all of the chip's state be accessible somehow from off chip.

One of the problems encountered in many multichip microprogrammed machines is that a great deal of interchip communication is required in their operation. Although the bandwidth of the machine can be made large by pipelining the operations, any operation which requires the full circle through the feedback loop of the state machine will require a great deal of time for its execution. In the OM2 system, hardware features have been included in both the data path and the controller to reduce the chip-to-chip communication as much as possible. As already mentioned, the loop counter circuitry was included on the controller chip, which reduced the loop operation time from 3 cycles to 1 cycle. Chapter 5 mentions the conditional ALU operations in the data path which can modify the actual ALU operation as a function of the flag bit. An example of the utility of this capability is provided by the multiply operation. When performing a multiply, the ALU should either add two numbers or just pass one of the numbers straight through, depending upon the state of a flag. One way to do this operation would be to send the
flag to the controller chip, and execute a conditional branch to one of two locations. One of the two appropriate ALU operations would be at each of the two microcode locations. However, it would take several cycles to perform each step of a multiply were this method used. Since in OM the ALU on the data chip has the capability of modifying its instruction as a function of the flag, it actually takes only one cycle to perform this part of the multiply step.

There are times when it would be convenient to communicate many bits between the controller and the data path in one cycle. For instance, when emulating the instruction set of a higher level machine, the data path can examine various fields in the instruction currently being emulated and calculate microcode branch locations. It is then necessary to load the μPC latch with the calculated branch location. To facilitate this loading, a 16-bit bus connects the two chips, and is referred to as the "literal bus". To economize on the data path's pin count, when this bus is not transferring literal data between the two chips it is used to load microcode into the data path chip. A large number of pads are required for the microcode and data path literal interconnections. There was insufficient space along the left edge of the chip for all of the pairs of pads required for this communication. Hence, some pads were placed along the top of the chip and others along the bottom, and connections between these pads and the buses were made by running vertical wires to the appropriate bus lines where they run between the two stacks.

The layout of the completed controller chip is shown in figure 13. A floor plan of the controller is given in figure 14, for use as an aid in studying the layout figure. Examples of the use of some of the controller's operations are given in the following section.

**Examples of Controller Operation**

This section will illustrate the operation and programming of the controller presented in the last section through the use of four programming examples: subroutine linkage, For-loops, Do-loops, and field dispatches. Refer to Table 2 for a tabulation of the controller's opcodes. It should be noted that the μPC operations are pipelined by one cycle so that if one particular microinstruction contains a controller jump opcode, the following microinstruction will also be executed before the jump actually occurs.

To call a subroutine, we would like to save the current value of the μPC on the μPC stack and load the μPC latch with the microcode address of the subroutine. When we have finished executing the subroutine code and wish to return, we just pop the return address off of the μPC
Fig. 13. The Layout of the Controller Chip

Fig. 14. Map of the Controller Chip
stack and load it into the μPC latch. To save the μPC value on the stack, "μPC" should be selected as the stack source and "PUSH" should be selected as the stack operation. As Table 2 shows, the condition must be true in order for the stack to push a value. Therefore, the condition selection should be "TRUE" to guarantee that the stack will save the return address. While we are saving the current μPC value on the stack, we must also load the μPC latch with the subroutine address. To do this, we select "MICROCODE" as the μPC source and put the subroutine address in the literal field of the microcode. Since we are not using the counter, the counter operation should be "NOP". For the return, we load the μPC latch with the return address by selecting "STACKTOP+1" as the μPC source and pop the stack by selecting "Pop" as the μPC stack operation. In order to guarantee that the stack pops the old value off the stack, we must make sure the condition is true by selecting "TRUE" as the condition selection.

Figure 15 illustrates the execution of subroutine linkages. Four "snapshots" of the microcode and μPC circuits are shown at the various steps as the execution proceeds. Snapshot (a) gives us a background for what is happening: The μPC is stepping through a segment of microcode, and is about to execute a CALL operation. The CALL operation contains a pointer to a sub-program located somewhere in the microcode memory. Snapshot (b) shows the state of the machine just after the CALL operation is executed. The μPC now points to microcode addresses inside the subroutine, while the return address to the main "program" is saved on the stack. Snapshot (c) shows that the μPC has advanced to the end of the subroutine, and the RETURN operation is about to be executed. The return address is popped off of the stack and loaded into the μPC latch, and program execution resumes where it left off in the main program, as shown in the last snapshot.

A For-loop should execute the same section of code many times. We can use the loop counter to store the number of times we have executed the code so that we know when we have finished the specified number of executions. Thus, when starting a For-loop, we should push the repetition number onto the loop counter stack. At the end of the loop we decrement the count, and if the result is not zero we should jump back to the start of the loop. If the decremented result is zero, we have finished execution of the For-loop, and we should pop the count off of the loop counter stack. Executing the For-loop in this manner requires that the end-of-loop command contain the address of the start of the loop. How then can we construct relocatable code containing For-loops? We can eliminate the need for the end-of-loop command to contain the loop's start address, by saving the start address on the μPC stack. The μPC latch would just have to be
loaded with the value contained at the top of the stack. Using this method of saving the loop address, the start-of-loop command becomes:

- \( \mu PC \) Source = \( \mu PC + 1 \)
- \( \mu PC \) Stack Source = \( \mu PC \)
- \( \mu PC \) Stack Operation = Push
- Condition = True
- Counter Operation = either Push Microcode or Push Literal

and the end-of-loop command becomes:

- \( \mu PC \) Source = Stacktop + 1
- \( \mu PC \) Stack Operation = Pop
- Condition = Count NOT EQ 0
- Counter Operation = True: decrement; False: Pop

The operation of For-loops is illustrated in figure 16. Again, four snapshots are shown which represent the state of the controller and microcode at various points in the execution of the loop. Snapshot (a) shows the state of the machine just prior to the execution of the FOR operation. When the FOR operation is executed, the value in the \( \mu PC \) latch is pushed onto the \( \mu PC \) stack, and the number of iterations specified by the FOR command is pushed onto the counter stack. The \( \mu PC \) continues advancing through the microcode. Snapshot (b) shows the state of the controller and microcode at some point in the middle of the FOR loop execution. When the end of the loop is reached, the value on the top of the counter stack is decremented. If the result is not zero, the new value is pushed onto the stack and the \( \mu PC \) latch is loaded with the value on the top of the \( \mu PC \) stack, as shown in snapshot (c). Notice that the value is not popped off the top of the \( \mu PC \) stack, because we will need the loop address again if the loop is not completed after executing one more time. When the result is zero, data is popped off the top of both stacks (to remove the loop address and the old count, which is now = 0) while the \( \mu PC \) value is just incremented, causing the controller to exit from the FOR loop, as shown in the last snapshot.

The Do-loop is similar to the For-loop, except that the code is repeatedly executed until a condition becomes true. That condition may be, for instance, when the data path flag becomes true. In this case, the condition selection in the end-of-loop command becomes "DATA PATH FLAG" instead of "COUNT NOT EQ 0". Also, since the counter is not being used, the counter operation in both the start-of-loop and the end-of-loop commands becomes "NOP".

Figure 17 shows some snapshots associated with the execution of a Do loop. By comparing figures 16 and 17, the similarities between FOR loops and DO loops can be observed. Basically,
(a) Before Execution of the "Call" Instruction.

(b) Just After Execution of the "Call" Instruction.

(c) Just Before Execution of the "Return" Instruction.

(d) After Execution of the "Return" Instruction.

Fig. 15. Illustration of Subroutine Linkage.
Fig. 16. Illustration of the Operation of "For" Loops.
(a) Before Execution of the "Do" Instruction.

(b) Execution of the "Do" Loop.

(c) Repeat Execution of the "Do" Loop While Condition is False.

(d) Exit "Do" Loop When Condition is True.

Fig. 17. Illustration of the Operation of the "Do" Loop.
Fig. 18. Illustration of the Operation of the "Dispatch" Instruction.
the only difference between these two types of loops is the decision of when to exit the loop. In a FOR loop a counter decides when the loop should be exited, while in the DO loop a flag, such as the flag from the Data Path, decides when the loop should be exited. Since the DO loop does not use the counter, the counter is not shown in the snapshots of figure 17.

When emulating the instruction set of a higher level machine, it is often convenient to do a multi-way branch. Suppose, for example, that the machine we are emulating has a 16-bit instruction word that contains a 4-bit opcode field and a 12-bit address field. In this case, we would have 16 code segments in the microcode, one for emulating each of the 16 possible opcodes of the higher level machine. We would like to be able to perform a 16-way branch, depending on the contents of the 4-bit opcode field, that would take us directly to the correct microcode segment, thus implementing the decode stage of instruction interpretation. We could use the ALU in the data path for calculating the microcode address for the proper segment, and load the $\mu$PC latch with the result of this calculation. This works especially nicely if the starting addresses of the segments are evenly spaced, because to calculate the branch address we merely multiply the 4-bit opcode by the segment length and add the displacement of the first segment. The multiplication is particularly easy to perform if the segment length is a power of 2, because then we just have to shift the 4-bit opcode value the appropriate number of places to the left.

A problem with the above method of field dispatching is that the microcode segments have to be evenly spaced in the microcode, preferably by a power of 2. In practice, segments are seldom of the same length. Even if they were of the same length, if one of the segments had to be modified, extensive corrections might have to be made all through the microcode. As an alternative, a dispatch table can be inserted into the microcode, which just contains a series of jump instructions to the appropriate microcode segments. If this is done, the 4-bit opcode value need only be shifted left once (because jump instructions are two microcode words long due to pipelining), added to the dispatch table displacement, and loaded into the $\mu$PC latch. To load the value into the $\mu$PC latch, the data path sends the result of the above calculation across the literal bus to the controller, and the controller selects a $\mu$PC source of "LITERAL".

Figure 18 illustrates the operation of the dispatch instruction. The controller jumps to a location in the dispatch table that is a function of one of the fields in the opcode. The dispatch table contains JUMP instructions to the various routines that perform the micro-instructions necessary to emulate each of the possible opcodes. The selection of the proper field in the opcode and the calculation of the dispatch table address are performed in the Data Path prior to the dispatch.
Some Reflections on the Classical Stored Program Machine

In the future, very large quantities of computing machinery may be placed on a single chip. Such chips will be easily and quickly designed, and rapidly implemented. This capability will present both a great opportunity, and a great challenge. How are we to organize and program such a wealth of hardware? Certainly not the way we do now.

Scaling of the technology to higher densities is producing effects which may be clarified by analogy with events in civil architecture. Decades ago, standard bricks, "two-by-fours", and standard plumbing were used as common basic building blocks. Nevertheless, architects and builders still explored a great range of architectural variation at the top level of the time: the building of an individual home. Today, due to the enormous complexities of large cities, many architects and planners have moved on to tackle the larger issues of city and regional planning. The basic building blocks have become the housing tract, the business district, and the freeway network. While we may regret the passing of an older style and its traditions, there is no turning back of the forces of change.

In present LSI, where we can put many circuits on a chip, we are like the earlier builder. While we no longer tend to explore and locally optimize at the circuit level (the level of bricks and two-by-fours), we still explore a great range of variation at the level of the individual computer system. In future VLSI, where we may put many processors on a chip, architects will, like the city planner, be more interested in how to interconnect, program, and control the flow of information among the components of the overall system. They will move on to explore a wider range of issues and alternatives at that level, rather than occupying themselves with the detailed internal structure, design, and coding of each individual stored program machine within a system. If systems are to work at all, they must at the least be understood at their highest level. These are some of the issues explored in chapter 8.
References


3. M. V. Wilkes, "The Best Way to Design an Automatic Calculating Machine", address to the Manchester University Computer Inaugural Conference, July 1951. The basic principles of microprogramming were first stated by Wilkes in this address.


Chapter 7: System Timing
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Sections:
The Third Dimension • • • Synchronous Systems • • • Clock Distribution • • • Clock Generation • • • Synchronization Failure • • • Self-Timed Systems • • • Signaling Conventions • • • Synchronous Elements • • • Asynchronous Elements • • • Arbitration

The Third Dimension

The successful design of large scale integrated systems requires careful management not only of the two-dimensional silicon area, but also of the operation of the system in the time dimension. Although time is physically different than the spatial dimensions, the general strategies already introduced for carrying the spatial design from conception to layout apply to system timing as well. These are the usual strategies for containing complexity: use of abstraction and structured design.

Much of the functional design of the spatial aspect of a system is done with the help of block diagrams, logic diagrams, circuit diagrams, and stick diagrams, in a metric-free topological domain. These representations are helpful because they allow designers to suppress detail, so that they can think about system behavior at a level of abstraction which is effective for the task at hand. One specific abstraction employed in these diagrammatic representations is the suppression of geometrical detail, while focussing on the topological structure of the circuit or system. Topology is sufficient to specify information flow between functional parts, so diagrammatic representations are a useful abstraction to the functional or logical structure of a system.

The third dimension, time, may also be regarded as having features analogous to geometry and topology. The definition of a sequential process -- whether represented by a program, flowchart, state diagram, or in plain English -- specifies only the ordering, or partial ordering, of the individual steps that compose it. Thus it is the metric-free "topological" concept of sequence, rather than the physical concept of a time metric, that is most useful for the functional design of a system.

As was pointed out in section 3-[Relating Different Levels of Abstraction], it is important that the levels of abstraction be related to each other and to physical concepts. The sequence domain is a self-
consistent abstraction that applies across several levels of system design -- programming, organization, logic. However, flowcharts and state diagrams do not say anything explicit about space, time, and other physical characteristics of a system, just as logic diagrams do not. The value of abstraction to the design process is that it permits one to defer certain bindings to physical form. The hazard is that one can become so isolated from physics and economics as to produce elegant schemes that are unworkable in practice. Thus, an important goal of any study of timing is to devise and explain methods by which sequence and time can be related.

At some level in the mechanization of a sequential process, one may no longer ignore the time metric. The electrical behavior of devices and wires is governed by physical laws which are expressed as partial differential equations in time. Devices and wires also take space, and their temporal behavior depends on these geometrical aspects of their construction. It is not generally possible, therefore, to separate the spatial and temporal aspects of system design. Failure to account for physical delays in the implementation of systems frequently results in unreliable operation, poor performance, or both.

Time is also important to people. We believe that the process of design starts with a conception of the functional operation of a system, together with a set of requirements -- or desires -- expressed in metrical units of space and time. These requirements are determined not by physics, but by human needs, expectations, and desires. The "interactive" text system that requires several seconds to respond to a keystroke is misnamed.

Unfortunately, the world is full of examples of digital systems which -- even when functionally correct -- have disappointed their designers and users as being unreliable or too slow. Why is it that so many systems have "timing problems", or fail to achieve performance objectives?

As is the case with the spatial dimensions, the design problems in the third dimension result not from a lack of possible forms, but rather from an overabundance. If one is to build a large scale integrated system with any hope of correct operation, it is necessary to restrict oneself to a consistent style of design. The canonical forms in the time dimension are signaling conventions which are adhered to throughout the system, and serve the function of establishing between all parts engaged in a communication an interval or sequence of intervals of time for this communication. If such a scheme is to be regarded as a discipline, it must be possible to state precisely the requirements that the signaling convention places on system interconnections and element timing.
Alternative disciplines of design in this dimension can be characterized by the way in which they connect sequence and time. Let us take as an example the synchronous discipline of design, which has been used in a form with a two-phase clock in the designs presented previously in this book. Here, sequence and time are connected by means of the clock. The term "synchronous" comes from Greek: syn, or sym = same, or together + chronos = time; and the discipline is well named, since it requires all parts of the system to operate together in time. Since synchronous systems are by far the best known and most widely used, we take them as the starting point for the body of this chapter.

Synchronous systems possess some serious limitations, which are made even worse as \( \lambda \) is scaled down, and as systems become larger. One problem is efficiency. It usually happens that most of the combinational paths are short, and the system clock period is determined by one or a few seldom used slow paths. One particularly difficult situation with slow paths occurs when synchronous signals must be driven off-chip. The time required to drive a signal off-chip is today a substantial fraction of the minimal clock period of about 100\( \tau \). As \( \lambda \) is scaled down, the \( \tau \)-relative delay off-chip gets larger, indeed may exceed the minimal clock period. So, it appears that synchronous communication across chip boundaries will become less and less attractive. Synchronous communication within a chip appears to be at least possible down to the fundamental limit of about 0.25\( \mu \) channel lengths, but it would be very difficult to manage a synchronous design of the number of parts implied by this scaling while achieving reasonable efficiency.

The same considerations of managing the design of very large integrated systems which provide a motivation for dividing a system into modular parts argue that the parts be independently timed. If the parts are each synchronous systems with independent clocks, information communicated from one part to another must be synchronized to the receiver's clock. Unfortunately, as we show in a later section, this synchronization cannot be accomplished with complete reliability. The reason for this problem is that synchronizing elements are bistable, and have a metastable or balanced condition that occurs under the conditions in which synchronizers must operate. As was discussed in Chapter 1, there is no bound for the time the bistable element may remain in this metastable condition. There are many methods to reduce the probability that such a fault would crash a system, but all cost time and so reduce efficiency.

The limitations imposed by the synchronous discipline suggest that other disciplines be tried. The final sections of this chapter are devoted to an outline of an alternative called self-timed logic.
This term refers to a discipline of digital system design in which the timing aspect of the design is confined to the interior of elements. Elements can be designed in a number of ways, for example, as synchronous systems with an internal clock which can be stopped synchronously and restarted asynchronously. This type of clock allows synchronous elements to communicate reliably, because their clocks are partly dependent, i.e. not independent. Elements may also be designed as speed-independent asynchronous circuits. In any case, the terminal behavior of a self-timed element must satisfy a sequence domain representation, which assures that correct sequential operation of a self-timed system is insensitive to element and wiring delays. There are no clocks or global time references in a self-timed system. Instead, initiation of a given computational step depends on completion signals produced by its sequential predecessors. Thus self-timed systems operate at a rate determined locally by element and wiring delays, a rate which tends to reflect typical rather than worst case delays.

The subject of self-timed logic has two principal facets, the design of elements and the design of systems of interconnected elements. Along the seam between these subjects are conventions for delay-insensitive signaling. This bifurcation of the discipline is deliberate. The design of elements is difficult because it is here that logic, physics, and timing come together. However, the element designer can work within a domain in which physical and logical scale are both restricted to be small enough to make the design manageable. The design of systems is difficult because of the combinatorics of scale. However, the system designer can work within a domain similar to that of a programmer, in which many of the details of the underlying physical system have been suppressed and replaced by an abstraction which is free of hidden rules -- namely, the sequence domain abstraction.
Synchronous Systems

In the synchronous discipline of design, sequence and time are connected through the use of a system-wide clock signal. The clock signal serves two purposes -- or one might say it serves two masters. The clock is a global sequence reference, and is also a global time reference. As a sequence reference, its transitions serve the logical purpose of defining successive instants at which system state changes may occur. As a time reference, the period or interval, either fixed or variable, between clock transitions serves the physical purpose of accounting for element and wiring delays in paths from the output to input of clocked elements.

The ability of the clock signal to serve two masters, logic and physics, has a certain compact elegance, and conforms to an obsolete tradition of parsimony in the use of active elements. However, the dual role of the clock binds the system sequencing and timing so closely that "timing" is the source of numerous difficulties in the design, maintenance, modification, and reliability of synchronous systems.

The logical model which synchronous systems resemble is the finite-state machine, a model that has been described in detail in section 3-[Finite-state Machines] and in chapter 6. As illustrated in Figure 1, any such system must satisfy a topological requirement that every closed signal path pass through a clocked storage element. Closed paths which do not pass through clocked storage elements are excluded as they may create non-deterministic behavior either through oscillation or through asynchronous latching. There are several important consequences of this topological constraint on the logical design. First, it assures deterministic behavior if the physical aspects of the design are also correct. Second, it relieves the designer of any requirement that the combinational logic be free of transients -- static or dynamic hazards -- on its outputs. The only dynamic characteristic of a combinational net that matters is its propagation delay time. Finally, the storage or history dependence of the system resides entirely within the clocked storage elements, a fact which simplifies the design process and often also the maintenance and testing of a system.

The clocked storage elements in a synchronous system may take any of a variety of forms, discussed below, depending on physical requirements such as speed, economy, or static operation. While these elements are distinguished as being the only recipients of clock signals, in practice there may be a number of timing signals derived as different phases or submultiple frequencies of the clock. In these cases it may be difficult to see the correspondence between the circuit and the
finite-state model. Circuits such as shift registers are finite-state machines, but already possess such a natural and regular structure that it would be pointless and awkward to describe their behavior with state diagrams. Control elements, such as the finite-state machine stoplight controller described in section 3-[Finite State Machines], are a case of imposed structure, in which the combinational logic and clocked storage elements can be patterned on the silicon in a form that mimics the usual block diagram of a finite-state machine.

Clocked storage elements may be clocked by different schemes, but all are binary storage devices. The sort of physical device that has the property of storing information -- also called memory, or history dependence -- are those which store energy, or are those such as film or punched cards in which energy is required to change some detectable condition of the medium.

For semiconductor integrated circuits the energy represented by charge stored on circuit capacitances is the only practical mechanism for storing information. Inductance plays the same role in superconducting circuits. MOS circuits employ this mechanism very directly in the dynamic register introduced in Chapter 3 and used in designs throughout this book. In the dynamic register illustrated in Figure 2, the output stored data follows the input as long as the enable input to the transfer gate is high. When the enable signal goes low, the charge stored on the node is very well isolated, and so maintains the same voltage.

Unfortunately, this is not the whole story. While the charge is very well isolated, it is not perfectly isolated. Charge escapes by two different mechanisms, which scale differently. The principal leakage path for 1978 MOS technology is the reverse leakage current of the drain junction of the pass transistor. The time constant of this decay is in the order of a few seconds at room temperature, but decreases exponentially with temperature to a millisecond or so at 70°C. This leakage path is a current per unit area, and because scaling down the circuit dimensions increases the capacitance per unit area due to decreased oxide thickness, the time constant of charge decay increases with reduced circuit dimensions. The scaling is largely masked by the exponential temperature dependance, however, so the time constant of junction leakage is reasonably regarded as constant for values of $\lambda$ over the recent past and future. Subthreshold currents are expected to become the limiting factor in holding charge on a node as soon as threshold voltages are reduced to much below 1 Volt. This effect of scaling down the dimensions of MOS circuits was discussed in section 1-[Effects of Scaling Down the Dimensions of MOS Circuits and Systems]. We refer to the time a node will reliably hold a high level as the refresh period.
Figure 1: Finite-state model
Figure 2: Dynamic register (1 bit)
The decay of charge on a dynamic node is no problem so long as the charge is sensed and
refreshed frequently enough, for example in every clock period as is done in the CM design. In
some cases this is not possible, and so-called static registers are used instead. An assortment of
static registers is shown in Figure 3. The general idea evident in the first two circuits is to
amplify and feed back the stored data so as to counteract the decay of charge. The first circuit
does this through a resistance which must be much larger than the effective $R_{on}$ of the pass
transistor, so that the storage node can be driven to the value desired. For voltages close to the
switching threshold this circuit amounts to a negative resistance termination to $V_{inv}$, where the
resistance is $-(\text{large } R)/(\text{voltage gain of the pair of inverters})$. This circuit can be used
advantageously as a termination for buses to assure static operation. The third diagram is a
storage circuit typical of bipolar families, in which low impedances make dynamic storage
unattractive, so that static storage is the rule rather than the exception.

These static circuits are logically equivalent to the dynamic storage circuit, except that the stored
information will remain indefinitely. The use of extra circuitry for each node to accomplish this
continuous refreshing is usually unnecessary, and it will be omitted in the following discussion
and figures with the understanding that it could in all cases be included where required.

As an aside, the reader may wonder about the statement above that capacitance provides the
mechanism for information storage. Is this true of the static and cross-coupled storage circuits as
well? Many references on switching circuits leave the impression that the existence of two
logically consistent stable states in these cross-coupled circuits is sufficient to assure that the
circuit will store a bit. Some references mention also that the circuit must have more than unity
gain around the loop, which is indeed a necessary but not sufficient condition. Consider the
consequence if the circuit capacitances were all taken to zero. The circuit would then be able to
respond instantly to an excitation, which means that a current pulse of arbitrarily short duration
could change the state of the circuit. If this arbitrarily small amount of energy could change the
state of the circuit, one could not reasonably expect the circuit to remain in either state. Without
capacitance, it cannot store information. This physical aspect of storage devices is discussed in
detail in Chapter 9.

With this background on storage elements, let us proceed to clocking schemes. As we pointed
out in the opening section of this chapter, it should be possible to state precisely the requirements
the timing form places on system interconnections and on element timing. For synchronous
systems, the requirement on interconnections is the topological requirement that all closed paths
pass through a clocked storage element. The requirements on element timing depend on the clocking scheme.

The storage devices described above, all logically equivalent to the dynamic register, may be used in a cheap, fast and risky clocking scheme illustrated in figure 4. We know of no example of this scheme being used in LSI circuits, but it was common in many of the early "transistorized" computers, with the gated, cross-coupled storage element shown in figure 3. This scheme might best be termed "narrow pulse clocking," because it requires that the clock pulse be narrow compared to the delay of the combinational logic. The present state information changes a short time, about $R_{on}C_{in}$, after the leading edge of the clock. The delay through the combinational logic must be greater than the clock width, or else the change in the present state information will propagate through the combinational logic to change the next state information before the trailing edge of the clock.

Once a clock period and width are established, the combinational logic must be designed to satisfy a two-sided bound on its delay time -- greater than the clock width and less than the clock period. As indicated above the clock width is also bounded on two sides -- below by the time required to transfer charge to the present state inputs of the combination logic, and above by the minimum delay of the combinational logic. The clock period also has a two-sided bound, unless static registers are used. The relations which must be satisfied are summarized in figure 4. They are relations which apply in a worst-case sense, in spite of variation in temperature, power supply voltage, aging, and manufacturing.

This "narrow pulse clocking" scheme was abandoned because of the difficulty of satisfying so many two-sided bounds simultaneously under so many conditions of variation. Also, the economies achieved due to the simplicity of the clocked elements were partly offset by necessity to "pad out" the delay in many of the combinational nets. This clocking scheme is quite feasible for certain LSI systems, since it is inherent in their manufacture and operation that most of the variables will track if the clock signal is generated on-chip.

Two-sided bounds on timing create enough difficulties in the design and maintenance of a system that it is generally worthwhile to use more logic to make the timing bounds one-sided. Elimination of the two-sided bound on clock width, or the complementary bound on combinational delay, requires the use of at least two clock phases. This minimum form occurs for much the same reason that ship canal locks require at least two watertight gates.
Figure 3: Static Registers
Figure 4: Cheap, fast, and Risky Clocking Scheme
The two-phase clocking scheme illustrated in figure 5 includes four sequentially repeated epochs. During \( q_1 \), previously stored information is applied to the present state inputs of the system's combinational logic. The \( q_1 \) signal must remain high long enough to charge the present state input nodes, a process which incurs some delay on the order of \( R \cdot C_{in} \), and which is called the delay time of the clocked storage element. Following this delay, if inputs are also available, the combinational logic starts setting up the outputs and next states, independent of when \( q_1 \) may transit from high to low. This epoch is analogous to the operation of a canal lock releasing a ship. The gate must open before the ship leaves, and the ship must clear the gate before it is again closed. If the lock master chooses to leave the gate open for a while after the ship leaves, it does not slow down the ship.

What the lock master must never, never do is to open both gates at once! The epoch labeled in figure 5 as \( t_{12} \) is an interval produced by the non-overlapping phases of the clock. By analogy with the narrow pulse clocking scheme, it is clear that overlap less than the minimum delay of the combinational logic is harmless to correct operation. However, as the minimum delay of the combinational logic is ordinarily legislated to be zero -- most people would agree it could not be less, and for circuits such as shift registers the delay does approach zero --, the overlap period \( t_{12} \) must be greater than zero. In practical cases, because \( t_{12} \) does not represent "dead time", but time during which the combinational logic is working, this time is made as short as is convenient but not necessarily as short as is possible.

In the epoch during which \( q_2 \) is high, the clocked element samples its input. The combinational outputs must be stable slightly before the trailing edge of \( q_2 \), an interval called the preset time of the clocked storage element. After this point in time, changes at the input of the pass transistor clocked by \( q_2 \) will not be fully passed to its output. Of course, \( q_2 \) must be wider than the preset time. This \( q_2 \) epoch is analogous to the entry of a ship into a lock. The ship cannot enter the lock until the gate is opened, and the gate should not be closed until the ship is completely inside.

Following \( q_2 \) is another period of non-overlap, \( t_{21} \), during which the system is idle. The minimum clock period for correct operation is the maximum combinational delay, plus the maximum delay time and preset time, plus \( t_{21} \). It is important therefore to make \( t_{21} \) as small as possible if one is designing for performance. As we show in the following section, \( t_{21} \) does serve a useful purpose of accommodating clock skew, a variation in the arrival time of the clock to
different clocked storage elements, and so in some systems $t_{21}$ can be reduced only as much as a clock skew allows.

The net result of this two-phase clocking scheme is that the clock period and its constituent epochs are, with static storage devices, bounded only below. A region of correct operation can always be found by making periods larger. With dynamic storage devices, the upper bound on the period is so large compared with the lower bound, at least for present values of threshold voltages, that the region of correct operation is always adequate. The complementary requirement on the propagation delay of the combinational logic is a simple upper bound.

Many variations on this basic scheme are found in different kinds of digital systems, more schemes than we could hope to describe individually. Many processors and storage systems are advantageously designed with more than two clock phases. In processors these phases -- usually four or eight, and sometimes a variable number -- delineate minor cycles which subdivide the major cycle. The general rules here are quite simple. No path may lead from a register or storage element output to an input clocked on the same phase. The maximum propagation delay time of a path from a register clocked on some phase, say $\varphi 3$, terminating at the input of a register clocked on some other phase, say $\varphi 6$, cannot exceed the minimum time between the leading edge of $\varphi 3$ to the trailing edge of $\varphi 6$, less delay and preset times. Multiple phases are also commonly used in systems which employ precharged pullup, and other charge transport techniques such as CCD storage. Magnetic bubbles are made to move in response to a rotating magnetic field, a two-phase clocking scheme when viewed as two orthogonal fields with sinusoidal oscillation $90^\circ$ out of phase with each other.

Some mention of variant forms with respect to inputs and outputs is also required. Inputs to synchronous systems must appear in synchrony with the clock, a requirement which is readily satisfied when the inputs are outputs of another system which shares the same clock. If an input does not assume its correct value until after the leading edge of $\varphi 1$, its worst-case delay relative to the leading edge of $\varphi 1$ is accounted for in the same way as the delay of the clocked storage element. The general procedure for checking compliance with timing bounds consists of marking nodes starting with clocked element outputs and system inputs with the latest time relative to the beginning of $\varphi 1$ that the signal will become stable. Programs to accomplish such checking are not trivial, as they may at first appear, because the program should account for different delays in different states and inputs. Otherwise, the large differences between delays for positive and negative transitions cannot be accounted for; nor can circumstances such as time of output

Figure 5: Two-phase Clocking
determination in simple AND and OR circuits be dealt with simply, because these times are data dependent. Unfortunately, this checking problem is about as difficult as exhaustive simulation.

The form of finite-state machine model of synchronous systems used in the descriptions above of clocking schemes is the transition output (Mealy) machine. It is more general than the state output (Moore) machine in that the outputs are functions both of the input and present state, while for the state output machine the output is a function only of the present state\(^9\). Transition output machines tend to be used in cascade arrangements in which economy is the principal design goal, while state output machines are characteristic of pipeline architectures in which high performance is sought.

Networks of transition output machines must be acyclic, the same requirement as for the combinational paths in a single machine. Combinational delays may accumulate on paths through many machines, so the general checking procedure described above must be used. Networks of state output machines may be connected cyclically, and checking is confined to communicating pairs of machines. Checking can be made entirely local to each finite-state machine if communication between machines is performed in pipeline fashion through clocked elements. The finite-state machine described in section 3-[Finite State Machines, (figure 14)] is of this sort.

In a two-phase clocking scheme, \(q_1\) and \(q_2\) are symmetrical in the sequence sense, and perhaps also in time. This symmetry suggests a reversal of roles between \(q_1\) and \(q_2\) is possible. For example, there is no reason in the finite-state machine structure shown in Figure 5 not to replace the simple amplifier or double inverter with combinational logic. The general structure allowed in a two-phase clocking scheme is any composition of basic elements consisting of registers followed by combinational logic, in which outputs of elements clocked by \(q_1\) drive only inputs to elements clocked by \(q_2\), and vice-versa. Please note that combinational delays must be checked across communicating pairs of machines. This scheme is similar to that used in the OM2 and its controller, and is well illustrated in Chapters 5 and 6.
Clock Distribution

Readers who have been designing systems with catalog parts are now invited to stand up and object: "What is with all of these clock phases? My systems use a single-phase clock." That is a good question.

The clock supplied externally to catalog parts such as registers, counters, shift registers, and microprocessors is most often a single-phase clock because this approach is convenient and the clock then uses only a single package pin. Internally, a two-phase clock or its functional equivalent is derived from the single-phase clock, either as part of each clocked element with the single-phase clock distributed through the chip, or once for the chip with the derived two-phase clock signals distributed as required.

Figure 6a shows a relationship between a single-phase clock and a two-phase clock derived from it, and figure 6b is a circuit which performs this function. While other conventions of relating single-phase and two-phase clocks are possible, in what follows this form will be taken as canonical. The single-phase clock is used in a trailing-edge triggering discipline, so-called because system state changes occur following the trailing edge of the clock pulse. This mode of operation has many advantages over leading-edge triggering. For example, because the state variables are stable during the clock pulse, one can perform logical operations between the clock and state variables in order to derive gated clock signals for selective register loading.

The preset and delay time relations to the derived two-phase clock can be translated, as shown in figure 6a, to be referred to the trailing edge of the single-phase clock. A study of this figure shows that, because of delays in deriving \( q_1 \) and \( q_2 \), preset times referred to the single-phase clock may be negative. Over a set of clocked storage elements the preset times will have maximum and minimum values, just as delay times will. The maximum preset time is called the setup time, and the minimum preset time, with its sign reversed, is called the hold time. In single phase clocking arrangements, the input to a clocked storage element should become stable by the setup time before the clock edge, and hold this value until at least the hold time after the clock edge. In the two phase clocking arrangements most often used in MOS LSI, the preset time is always positive since it is referenced to the trailing edge of \( q_2 \), so only its maximum value or setup time is critical to correct circuit operation.

It is an essential feature of clocking schemes with a one-sided bound that there is a critical period
Figure 6a: Relation between a single phase clock and a two-phase clock

Figure 6b: Logic Circuit to derive two-phase non-overlapping clock from a single-phase clock
during which the clocked storage element is actually storing two bits of information. For either single- or two-phase clocking this critical period begins at the preset time and ends at the delay time. One of the stored bits is the input value presented before the preset time; the other is the bit presented at the clocked storage element output. This critical period provides a built-in tolerance for clock skew. This term refers to a variation in the effective arrival time of the clock at different clocked elements. These variations may be due to a combination of several effects: different threshold voltages, signal propagation delays on wires, or variation in element delays such as is found when gated clock signals are used to control register loading.

Clock skew even within a chip can be a problem. As was pointed out in section 2- Electrical Parameters, propagation of signals on poly lines, as they are of fairly high resistance, is a diffusion process for which the delays are not negligible. Given the distances and need for short delay in clock distribution, use of the diffusion layer for carrying the clock more than short distances is not recommended either. For example, the delay in a line 6μ wide and 6 mm. long is calculated by the method presented in section 1- Delays in Another Form of Logic Circuits, and using the typical 1978 MOS electrical parameters given in Table 1, section 2- Electrical Parameters, to be about 100 nsec. in 70Ω/□ poly and about 30 nsec. in diffusion. Propagation over the same distance on a 9μ metal line requires only 0.1 nsec. Other sources of clock skew are little problem. Threshold voltages of elements built as monolithic companions tend to track. Skew in signals produced by gating a clock is present but consistent, and can be controlled by a method presented later in this section.

Clock skew is a much more serious and common problem in systems built from "families" of SSI and MSI circuits. The origin of the problem is not entirely clock distribution, but the manufacturers' inattention to preset and delay time bounds. It is easy to find examples in popular families of catalog parts certain pairs of clocked elements in which the hold time (minimum negative preset time) specified for one part is greater than the minimum delay time of another part. Although the system may work if delays are "typical," some fraction of a production run can be expected not to work, or worse, may fail intermittently in service. Designers of families of circuits intended and advertised to work together should strive to make all preset and delay time characteristics identical within a family, as variations from one part to another decrease the tolerance of the system to clock skew.

Let us return now to the typical MOS integrated system with a two-phase clock distributed as the signals q1 and q2. These are generally "popular" signals, second only to VDD and ground, and
the capacitance accumulated in their distribution and gate connections is considerable. It is possible, of course, to drive \( p_1 \) and \( p_2 \) directly onto the chip from an external driver. Systems operating at the extremes of performance of the MOS technology benefit from this approach. The external clock driver can be made to switch fast, and to a higher voltage than VDD, so that transmission gate outputs can then transit all the way to VDD. Because the saturation current varies as \( (V_{gs}' - V_{th})^2 \), there is a large speed advantage on the chip of the higher voltage clock drive.

The performance benefits from driving \( p_1 \) and \( p_2 \) onto the chip from fast off-chip drivers are achieved at the expense of external components. Although on-chip clock drivers may limit the performance of a system slightly, they are generally used in the interests of economy. An integrated system which is large by today's standards may have a clock load on the order of \( 10^4C_g \), where \( C_g \) is the gate capacitance of a minimum dimension transistor. The techniques developed in section 1 [Driving Large Capacitive Loads] for driving large capacitive loads are applicable to clock driving. The total delay in an exponential driving structure with the optimum fanout of e, starting from a minimum energy signal to drive \( 10^4C_g \), is only \( 25\tau \). However, the last stage of the driver would be impractically large; the gate width would be \( 2'10^4\lambda/e \). The tradeoffs between performance and area may dictate a large fanout for the final stage of the clock driver, say about 40, with smaller fanouts for the drivers preceding it. Most of the delay in this driver structure would be in the final stage.

A clock driver of this sort is illustrated schematically in figure 7a, with size ratios indicated by the delay times. Waveforms expected for this circuit are shown in figure 7b. The clock is assumed to originate from somewhere off the chip. There is no reason to use two pins for this function where one would serve, so the canonical scheme for production of a two-phase clock from a single-phase clock is used. The capacitance of the pin and package are so large compared to \( C_g \) that there is no point is starting from a minimum energy signal. The structure shown presents a load of about \( 30C_g \). If the clock source were an on-chip clock generator, more stages would be used. The slight asymmetry between \( p_1 \) and \( p_2 \) is of course due to the inverter at the input.

The reader should take careful note of an important characteristic of this clock driver; namely, that non-overlap of the clock phases is assured independent of clock loading. This desirable characteristic is the reason that it is the clock phases, rather than the NOR gate outputs, that are fed back in a cross-coupled fashion. The non-overlap periods can be reduced somewhat at the expense of silicon area by cross-coupling at every stage. This same circuit trick can be extended
Figure 7a: Clock Driver

Figure 7b: Clock Waveforms
as shown to assure non-overlap independent of clock loading for gated clocks as well. The techniques used in these clock driver circuits to assure the existence of the non-overlap period independent of the effects of loading on element timing are much in the spirit of those techniques used in speed-independent and self-timed systems, in that the circuit adapts its temporal behavior to conform to a sequencing constraint.

A question is sometimes raised as to whether clocks for different sections of a chip should be buffered separately, as is often done at a circuit board level for systems built from collections of circuit boards. On a physical basis, the answer is no. It is best for two reasons that the clock lines be common throughout the chip. First, this approach minimizes clock skew. Second, since one must pay the same transistor area whether the driver is lumped or distributed, optimum driver design locates the stray capacitance of the clock distribution wires where the largest signal energy is available, which is after the final driver. There may be organizational reasons for distributing the final driver stage to locations close to sections being driven when some logical operations are performed on the clock signal.
Clock Generation

Where great clock accuracy is required, or where a system clock is distributed to many circuits, the clock generation task is external to the chip and is not the direct concern of its designers. The clock will ordinarily originate from an electronic oscillator circuit, the period of which is controlled by a crystal or some resonant network. Process variation in integrated circuit fabrication does not allow accurate resonant networks to be fabricated by usual means, but it is perfectly feasible, indeed essential for self-contained VLSI systems, to generate clock signals on the chip. It is best in approaching this subject to forget about electronic oscillator circuits, and instead to take a more basic approach originating with an understanding of what clocks are for.

As we have mentioned before -- and this is a principle that bears repeating on every opportunity --, the role of the clock in a synchronous system is to connect sequence and time. The interval between clock transitions, whether these transitions are on one or distributed over several wires, must be such as to permit enough time for the activities planned for that interval. When viewed in this way, a clock is more like a set of timers than like an oscillator. A model of the temporal behavior of the systems being clocked is built into the clock generator in the choice of times for the various timers.

The easiest way to build these timers is as chains of inverters. The propagation delay time of such a chain will of course vary with \( \tau \), according to the way in which the fabrication process, aging, temperature, and power voltage affect \( \tau \). However, these variations only make the inverter chain a better model of the system being clocked than a fixed timer would be, since on the same piece of silicon these variable factors are nearly the same for the clock and for the system.

It is helpful to distinguish between the two kinds of timers shown in figure 8. The first is a *symmetrical delay*, so-called because the propagation delay for positive and negative transitions at the input is about the same. The second is a logic network designed to produce as asymmetrical a delay as possible. A negative transition propagates through the delay in about \( 5\tau \) independent of length. A complementary form of *asymmetrical delay* is also possible, but to simplify the figures and symbology in what follows, we shall use only the form in which a low input resets the delay and a positive transition propagates slowly. The symbols shown in the figure allow for *taps* at various points along the delay.

Clocks which employ these delays as timers are all elaborations of the *ring oscillator* circuit shown
Symmetrical delay, delay $\approx 5n\tau$

Asymmetrical delay,
positive transition delay $\approx 5n\tau$
negative transition delay $\approx 5\tau$

Symbol for a delay with taps

Figure 8: Delays
Figure 9a: Ring Oscillator

Figure 9b: Symmetrical Clock

Figure 9c: Universal Clock
in figure 9a. Rings of an odd number of inversions have no stable condition, and will oscillate with a period which is some multiple of the delay time twice around the ring. The oscillation of the largest period will eventually predominate following bringing power on, but the erratic clock signals produced during power-up could leave the system in a peculiar state. It is much better to produce an initialization signal which is held high during power-up, and to use it to initialize the state and the clock. In the modification of the ring oscillator shown in figure 9b, clock signals are suppressed during initialization, and will start immediately following the negative transition of the INITIALIZE signal. This circuit produces a symmetrical single-phase clock which can be converted to a two-phase clock by the circuit shown in figure 7a.

Although the clock circuit shown in figure 9b would be adequate for many applications, many elaborations can be included which are shown together in figure 9c. The width of the single-phase clock pulse, related to the $q2$ period, is determined by one asymmetrical delay, while the interval between clock pulses, related to the $q1$ period, is determined by an independent asymmetrical delay. Another feature of this universal clock is that it allows the system being clocked to select between a variety of periods, which can be changed on a cycle-by-cycle basis according to the combinational delay of the operation performed on that cycle. In order to visualize how this works, note that following the trailing edge (negative transition) of the clock signal, any high input to the 5-input NOR circuit has the effect of preventing the occurrence of the next clock pulse. The usual default case is with the LONG', MED', and RUN high, and INITIALIZE low, resulting in a short period determined by the first tap on the period delay. If a decoding of the state indicates that a longer period is required for that cycle, the MED' or LONG' lines must be driven low before the short default period is elapsed. If for example the LONG' line were low, the period before the next clock pulse is stretched to that determined by the full delay. Of course, this scheme may be generalized to any number of delay taps. Signals such as MED' or LONG' can be derived from function coding of the combinational sections whose modeled delay they match, or as microcode bits.

The RUN line is a bus intended to generalize this cycle stretching feature so that any part of the system being clocked may stop the clock synchronously, and then permit it to restart asynchronously. If this cycle is ever to be stretched to more than the refresh time, static storage elements must be used (at least for the $q1$ part of the cycle for two-phase clocking). This technique of control over the clock is the basic mechanism exploited later in this chapter to allow asynchronous communication between synchronous systems.
Synchronization Failure

Jean Buridan (1295 - ? after 1366), a 14th century French philosopher often cited as a precursor of Isaac Newton for his priority in giving a technical definition of kinematic terms such as inertia and force, posed in his commentary on Aristotle's *De Caelo* a paradox; that a dog could starve if placed midway between two equal amounts of food. The unfortunate creature placed in this position would be equally attracted to each source of food, and in a position of equilibrium. One 20th century explanation of this paradox, if indeed it is a paradox at all, is that the structure consisting of the dog and the two sources of food is and behaves just as any other structure which can store a bit of information.

The analysis of the electrical behavior of cross-coupled circuits presented in section 1-[Properties of Cross-Coupled Circuits], and developed in physical terms in section 9-[Energetics of the flip-flop], applies also to the situation which Buridan described. The equilibrium condition either for the dog or for the cross-coupled circuit is unstable, as any displacement from equilibrium brings about forces which tend to destroy rather than restore the equilibrium condition. An unstable equilibrium of this sort is called a metastable condition. Buridan was correct in believing that the dog could starve, as it is characteristic of a metastable condition that it may persist indefinitely. A functional definition of metastability applied to cross-coupled circuits is the occurrence under undriven conditions of an output voltage in a range around $V_{\text{inv}}$ which cannot reliably be interpreted as either high or low.

A bistable element in a self-contained synchronous system never has the opportunity to reach a metastable condition, since satisfaction of the timing constraints assures that the output is driven to a voltage outside of the metastable range. But is any system really self-contained? A system such as a microprocessor may be entirely synchronous internally, but cannot extend this synchrony indefinitely to encompass all of the external world with which it may interact. If asynchronous signals of external origin are allowed to enter a synchronous system as ordinary inputs, the timing constraints required to assure correct operation cannot be satisfied, since there is no known relationship between the timing of the asynchronous inputs and the clock.

Figure 10 illustrates in a small fragment of a larger synchronous system the consequence of ignoring synchronization altogether. Even if one employs a model of the clocked storage elements as having perfectly discrete outputs, the unequal delay in the paths from the asynchronous input $X$ to the clocked storage elements allows the inputs to the clocked storage...
Figure 10: Fragment of a logic circuit illustrating the consequence of ignoring synchronization.

Figure 11: Fragment of a logic circuit illustrating a synchronizer which would work perfectly if perfectly discrete bistable devices were possible.
elements during state A to represent an illegal successor state for a period following a transition of the asynchronous input. If the clock happens to capture the inputs during this transitory period, one of the illegal state transitions shown in dashed lines on the state diagram will result.

So, a slightly smarter thing to do is to assure that only one clocked storage element is affected by a given asynchronous input. A clocked storage element that is used in this way is called a synchronizer, since it is intended to produce an output signal which is in synchrony with the clock. Figure 11 shows a redesigned version of the previous fragment of a circuit modified in its state coding so that the asynchronous input X affects the input to only one clocked storage element. If it were only possible to build perfectly discrete bistable devices -- indeed, if perfect discreteness exists in nature, for which, see chapter 9 --, this scheme would be perfectly reliable. Unfortunately, there is some probability of synchronizer failure, since a transition of the asynchronous input in certain times relative to the clock will leave the synchronizer in a metastable condition, and the time required for the clocked storage element to get out of a metastable condition is unbounded. During the period in which the synchronizer output remains in a metastable condition, the logic cannot discriminate between states B and C, and if the condition persists for too long, an illegal or incorrect successor state can result.

It is part of the interesting history of the subject of synchronization that even long after the necessity to synchronize asynchronous inputs was recognized as a standard part of good engineering practice, the faith of logic designers in the discreteness of the outputs of clocked storage elements was so great that the very existence of synchronization failure was widely denied. Another curious aspect of the sociology of the problem is the many schemes proposed to "solve" the problem, but which only move it to another location in a system or reduce its probability. Synchronization failure was discovered independently by numerous researchers, designers, and engineers in the 1960's, some of whom published reports of their analyses or observations\(^1\,2\,3\,4\,5\,6\). The work done at the Computer Systems Laboratory of Washington University by Thomas J. Chaney and Charles E. Molnar\(^7\), and by Marc Hurtado\(^8\) has provided convincing demonstrations of the existence and fundamental nature of the problem.

It is fairly easy to estimate the probability of a synchronization failure with a simple mathematical model. If one observes that a bistable device is in a metastable condition at some time \( t \), what is the probability that it will have left this condition by time \( t + \delta \), in the limit as \( \delta \) approaches zero? The only answer to this question that seems reasonable is that the probability is proportional to \( \delta \); let it be \( \rho \delta \). This assumption produces a simple model in which the exit from a metastable

\[ \text{[Ch7. System Timing. by C. L. Seitz: Sect.2]} \]
condition is a Poisson process of rate $\rho$, and the probability that a clocked storage element will remain in the metastable condition, once in it, for a period $D$ or longer is $e^{-\rho D}$. The prediction of this simple model has been verified experimentally and is consistent with analyses based on circuit models\textsuperscript{8}, including the analysis presented in section 1.4 (Properties of Cross-Coupled Circuits). The parameter $\rho$ depends on circuit characteristics. A dynamic storage element is not an acceptable synchronizer, as the time evolution of its output in undriven conditions is possibly even toward the metastable region (See section 9.7 (Energetics of the Flip-flop)). One can identify $\rho$ in the analysis in section 1.4 (Properties of Cross-Coupled Circuits) with $1/\tau_0$. The time evolution of the voltage output of the cross-coupled circuit has the effect of transforming a uniform probability distribution of initial conditions to an exponential or Poisson distribution of exit events. For ratio logic with a ratio of $r$, $\tau_0$ is about the pair delay, $(r+1)\tau$.

In order to estimate the probability of a fault in a synchronous system due to the non-vanishing probability that a synchronization will take longer than some bounded time, one must also calculate the probability that a synchronization event will put the synchronizer into a metastable condition. For most synchronizations, the asynchronous level to be synchronized will transit sufficiently far away from the time at which it is sampled that the clocked storage element will be overdriven in the usual way. Only over a rather narrow time aperture, denoted here as $\Delta$, does the occurrence of a transition result in the synchronizer taking more than the usual delay time of the clocked storage element. The boundaries of this aperture are not sharp, but may be treated as such, so that for a particular frequency of transitions of the asynchronous signal, $f$, the probability that a metastable condition will be produced in a single synchronization event is $f\Delta$. One may take this relation as the definition of $\Delta$.

The overall probability of a system failure at each synchronization event, $f\Delta e^{-\rho D}$, depends on $\rho$ and $\Delta$, which are parameters of the clocked storage element used as a synchronizer; on $D$, which is a parameter of the synchronous system in which the synchronizer is used; and on $f$, which is a parameter of the asynchronous input signal. $D$ is the time allowed in the synchronous system for the decay of the probability of metastability, and is effectively like a delay. It corresponds to the excess delay allowed from clocked storage element outputs to inputs (see figure 5), and even for zero combinational delay cannot exceed the clock period less delay and preset times.

In order to get some feeling for the failure rates involved, consider a synchronous processor which is accepting data from, or sending data to, a disc storage unit at a 1MHz rate. An asynchronous signal alerts the processor to the presence of, or need for, a new data item, but the processor is
able to clear this signal synchronously. We can assume that $r$ is about $1/5\tau$ for ratio logic with $r=4$. If almost all of a fairly short $100\tau$ clock period were available for the decay of the probability of metastability, $D$ would be about $80\tau$. It is interesting that when $D$ is expressed as a multiple of $\tau$, the exponent in the formula is then independent of $\tau$, and so is independent of scaling circuit dimensions. The scaled down version of this system would allow less time for the decay of the probability of metastability, but the synchronizer would exhibit a proportionately higher metastable exit rate.

We are not aware of any experimental determinations of $\Delta$ for nMOS circuits, but by analogy with experiments performed with several bipolar circuit families, we believe that $\Delta$ is a small fraction of $\tau$, say about $\tau/10$. This estimate agrees with the notion that $\Delta$ is time aperture corresponding to that time required for a signal to transit through a small voltage range around the switching threshold, a time which is proportional to $\tau$. For present values of $\tau$, $\Delta$ is then approximately 30 picoseconds, and the probability of a system failure for a single synchronization event would be about $(10^6)(30\cdot10^{-12})e^{-16}$ or about $3\cdot10^{-12}$. So, about one in each $3\cdot10^{11}$ items transferred across this interface would be in some fashion mistreated.

Failure rate depends on the frequency at which the system samples asynchronous inputs. This frequency cannot be greater than the clock frequency, and as is clear from a careful study of figure 11, this frequency may be as low as the frequency of the states whose choice of successor depends on the asynchronous input. This figure is intended to suggest how synchronizers can be sheltered from needless synchronization events, as the synchronizer is here sheltered by ANDing the asynchronous input with a signal which indicates that the system is in state A. Synchronizers that are used directly on asynchronous input signals and whose outputs enter PLA or ROM structures may cause failures even when the system is in a state in which the successor does not depend on the asynchronous input.

For the example above, the synchronous processor must sample every transition of the asynchronous input in order to transfer every data item. If this processor were engaged in transferring data at a $1\text{ MHz}$ rate about one third of the time, synchronization failures would occur at a (Poisson) rate of once each $10^6$ seconds, or about every 10 days. It is worth noting that the exponential relation in $\rho D$ makes the failure rate remarkably sensitive to $\rho D$. This dependence may be particularly noticeable if the clock signal originates off-chip so that $\rho D$ depends on $\tau$. A chip with a slightly larger than typical $\tau$ may exhibit a drastically higher than typical failure rate.
The probabilistic character of synchronization failures makes them exceedingly difficult to trace. Designers of synchronous systems who wish to avoid the curses and plagues that are the just reward for those that build secret flaws into human tools and enterprises should cultivate a rational conservatism toward this problem. The worst-case failure rate for a design should be calculated. If the failure rate is higher than some criterion, it can be reduced by techniques which increase D. Use of cascaded synchronizers is one technique for increasing D which does not require increasing the system clock period. Criteria for acceptable failure rates depend on many of the same economic and social factors that influence other aspects of system reliability.

One conservative failure rate criterion that can be supported by a physical argument is that the rate of synchronization failures should be on the order of the rate at which the bistable synchronizer will change state due to the random thermal motions of the electrons. This rate is shown in section 9 [Thermal Limit] to be \((1/\tau)e^{-(E_{sw}/kT)}\). If \(s\) is the frequency of synchronization events -- often \(s = f\) --, this physical criterion is:

\[
 sf\Delta e^{-\rho D} = sf\Delta e^{-D/(r+1)\tau} < (1/\tau)e^{-(E_{sw}/kT)}.
\]

Solving for D yields:

\[
 D > \frac{E_{sw}/kT + \ln(sf\Delta\tau)(r+1)\tau}{(r+1)\tau}.
\]

If one takes \(\Delta\) as about \(\tau/10\) and \(s\) and \(f\) in the order of 1/100\(\tau\), it is clear that the second term in the brackets can be ignored. The switching energy for today's circuits is about \(10^8kT\), and bistable devices are consequently so reliable that the time required to allow the probability of metastability to decay to achieve the same reliability is very large -- about \(5'10^8\tau\), or 0.15 seconds! However, this criterion scales in a remarkable way. The ratio \(D/\tau\), which represents the number of transit times the criterion provides for the metastable exit, scales with the switching energy, which goes down as \(\alpha^3\). This scaling should not be interpreted as meaning that smaller devices have a higher probability of metastable exit per transit time. Rather, smaller transistors result in less reliable storage devices, which makes it possible to lower one's standards. Ultimately small transistors with channel lengths of about 0.25\(\mu\) would allow circuits with a switching energy of about \(10^4kT\). Because of the significant subthreshold currents at the low threshold voltages implied by this scaling, CMOS circuits with \(r = 1\) would have to be used. At these minimum dimensions, this criterion implies \(D > 2'10^4\tau\), and taking \(\tau = 0.02\) nanoseconds, \(D > 400\) nanoseconds. So, at this ultimate limit of MOS technology, one cannot disregard synchronization failure, but one would not expect it to limit designs for synchronization rates up to 1 MHz, or so. Since this criterion represents the most conservative position that can still be rationally defended on physical grounds, it is known as the Mead Criterion.
Self-timed Systems

The operation of a synchronous system is reminiscent of soldiers marching to the commands of a drill sergeant. The temporal control over a collective activity is centralized in a single authority, and the soldiers respond to known commands that are synchronous with the marching cadence. Lockstep control results in a particularly simple form of organized behavior which people often associate with the relentless efficiency of machines. However, lockstep control is certainly not the only way to coordinate the collective activity of many participants, nor is it efficient unless the tasks of the participants are very well matched. Self-timed systems are patterned on quite a different image of organized activity, one in which the temporal control is delegated to the participants. If one were to try to construct a mental image of self-timed behavior, it would be one in which the airplane could not leave until after all the passengers scheduled for the flight had gotten on board. One tries to assure that everything occurs in proper sequence, but nothing ever has to occur at a particular time.

In a self-timed system, the metrical notion of time is confined strictly to the interior of parts called elements. Time and sequence are related inside these elements in such a way that the terminal behavior of an element obeys a set of sequence domain relations on the occurrence of signal transitions. The only kind of relations possible are \( x \leq y \), \( x \) precedes \( y \), for which one might also say either that \( x \) is before \( y \) or that \( y \) is after \( x \), and \( x \# y \), \( x \) is concurrent with (not ordered with) \( y \). The relation \( \leq \) is a partial ordering on the set of occurrences of signal transitions; it is reflexive, \( x \leq x \); antisymmetric, \( x \leq y \cap y \leq x \Rightarrow x = y \); and transitive, \( x \leq y \cap y \leq z \Rightarrow x \leq z \). The relation \( x = y \) in the definition above means that \( x \) and \( y \) are identical, not simultaneous. The notion of simultaneity is specifically regarded as meaningless and disallowed by the antisymmetric property of \( \leq \).

According to the special theory of relativity, it would not be possible to have any relations between occurrences of signal transitions at different points in space, as the order might be interpreted differently for observers in different locations. The problem with the relativistic environment is that one knows so little about time that systems are either impossible to make, or more complicated than is justified by the actual situation. Over volumes that are sufficiently small, volumes that are here called equipotential regions, one is justified in applying the approximations that (1) a signal is identical at all points on a wire, and (2) a relation which holds
anywhere in the region holds everywhere in the region. Of course, an element must be contained entirely within an equipotential region or else the set of relations which define its terminal behavior would be meaningless. A system may contain many equipotential regions. Each of them is conformal to a point in a relativistic space, and each keeps time to itself.

In a medium such as free space, events that originate at one point in a certain order are observed to occur in the same order at any other point. However, in a medium in which signals are carried on wires which may follow different routes from one equipotential region to another, order is not in general preserved. Some self-timed systems may use a multi-wire structure called a bundle, which can convey a set of signals from one equipotential region or element to another while preserving specified ordering relations. For example, a two-wire bundle with inputs (a,b) and outputs (c,d) may satisfy the constraint that \( a \leq b \Rightarrow c \leq d \), but would not also be required to satisfy \( b \leq a \Rightarrow d \leq c \), as the pair of requirements would imply that the bundle satisfy a two-sided timing bound.

The time required for an electromagnetic wave to traverse a large chip is less than 0.1 nsec, and for 1978 MOS circuits the signal transition times on-chip are greater than about 0.3 nsec. Accordingly, the equipotential environment is very well approximated on a single chip. For typical transition times on pins of many nanoseconds, a circuit board up to about a foot square is a good approximation of an equipotential region. However, communication delays in pin driving structures are so long compared with internal transition times that two chips may not be regarded as being in the same equipotential region.

<remainder of this section in preparation>

discussion of the formation of systems from elements based on the recursive definition: A self-timed system is either (1) a self-timed element, or (2) a legal interconnection of self-timed systems. Correctness proofs. Fundamental motivation for self-timed systems: that correctly functioning elements assure a correctly functioning system.
Signaling Conventions

\textit{in preparation}

two-cycle (transition) and four-cycle (Muller) signaling -- data validity -- the MOSbus convention -- ternary and data-driven signaling.

Synchronous Elements

\textit{in preparation}

illustration of the use of the stoppable clock to implement self-timed elements as synchronous systems.

Asynchronous Elements

\textit{in preparation}

self-timed elements implemented without clocks -- delay modeled timing -- speed-independent timing -- inherent error detection and correction.

Arbitration

\textit{in preparation}

asynchronous arbiters -- application to shared self-timed systems.

Acknowledgements

\textit{in preparation}

A. W. Holt, W. A. Clark, C. E. Molnar, Carver, Lynn, Ivan, Sproull, Dick, Wayne.
References


Chapter 8: Highly Concurrent Systems
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Sections:

Introduction - - Communication and Concurrency in Conventional Computers - - Algorithms for VLSI Processor Arrays - - Hierarchically Organized Machines - - Highly Concurrent Systems with Global Communications - - Challenges for the Future

How can the properties of VLSI be exploited to build computational structures? Our discussion to this point has concentrated primarily on principles for structuring circuits and wires on the chip rather than on the application of VLSI to solve interesting computational problems. Although the OM example described in chapters 5 and 6 shows an elegant use of the structuring principles in the design of a conventional processor, we are left with an intriguing question: Does VLSI offer more than inexpensive implementations of conventional computers?

This chapter answers the question with a resounding YES! Because processing elements and memory elements can both be easily implemented in VLSI, we are encouraged to find structures that use a great deal of concurrency—a large number of calculations occurring at the same time. Although we can clearly design VLSI structures that have many sites at which processing is performed, how are these structures to be applied? Some applications may require different sorts of concurrent processing than others. Are there any principles or theories that will guide us in the design of highly concurrent systems? (For an excellent introduction to the promises and problems of VLSI and concurrency, see reference [1.] ) Unfortunately, we lack experience in designing systems of this sort. As a consequence, this chapter can offer no complete designs which have been applied in real system applications. Instead, we offer several glimpses of the possibilities available with VLSI, and of its limitations.

The chapter is organized into four quite separate sections; although they are designed to be read sequentially, they may also be read concurrently! The first section reviews the problems that conventional computer designs present when implemented in VLSI, and summarizes efforts to achieve concurrency in general-purpose computers. Section 2 takes up a particular sort of concurrent organization—the array of identical processors—and shows its application to matrix
arithmetic. Section 3 examines hierarchically-organized machines--in this case, machines structured as a binary tree--and demonstrates how they can be programmed to perform several tasks. Finally, section 4 presents a nascent theory of planar computational structures. It links the topological and electrical properties of VLSI elements to the structure of computations.

1. Communication and Concurrency in Conventional Computers

The architectures of conventional computers suffer from two difficulties that we must avoid when designing VLSI computational structures. First, a processor is separated from its memory by long communication paths such as buses. These buses are long enough to slow substantially the transmission of information between a processor and memory. Second, the "von Neumann machine" provides only a single processor that sequentially fetches and executes instructions--it offers few opportunities for concurrent processing activity. In this section, we survey some of the attempts to reduce communication costs and to use several processors concurrently. Although designs using a great deal of concurrency have been cumbersome to implement in the past, VLSI makes these designs considerably more attractive because of the ease with which memory and processing elements can be placed in close proximity.

Human organizations, like computer organizations, suffer if communication costs are high or if concurrent processing cannot be exploited. In fact, a human brings to an organization what VLSI brings to a circuit: both combine processing and memory effortlessly! Analogies with human structures will help to suggest the kinds of behavior we might achieve in computational structures.

Humans struggle to reduce communication costs, because the cost is often measured in large quantities of time. Consider a student assigned to write a research paper, requiring the use of a large library. Each time he needs to consult a book, he could make a trip to the library, climb into the stacks to retrieve the book, read a few relevant paragraphs, and replace the book. He now heads home to write the sentence that depends on the information he acquired. Libraries and people both recognize the inefficiency of this approach, and allow students to borrow books. The student will take several dozen books home, and store them on a short shelf, handy to his desk. Now the communication cost required to find information is reduced, provided the item lies within the group of books he has selected. If the student finds it difficult to select a small
number of books that meet his needs, he may move his work to a carrel in the library, again in order to reduce communication costs with the large library "memory." The human strives to keep his information supply close to his processing task.

Concurrency is widely exhibited in human organizations. Henry Ford introduced the production line as a way to exploit concurrency in a well-understood manufacturing process. This is a particularly simple structure, in which information and goods flow rigidly along the production line. A more prevalent, general-purpose approach to concurrency in organizations is the hierarchy: the president of a company supervises several subordinates, each of whom in turn supervises a like number of sub-subordinates, and so forth until we reach the lowest level workers.

Two goals of the hierarchy are to keep everyone about equally busy, and to allow adequate information flow in the organization. A supervisor must generate enough commands to keep several subordinates busy—otherwise it would not be possible to build large organizations at all. In addition, each subordinate requires a certain amount of attention from the supervisor. These requirements limit the number of subordinates who can be assigned to a single supervisor—ten underlings can run the most diligent supervisor ragged. Supervisors gather information to make decisions by querying their subordinates. In a badly organized hierarchy, supervisors may confer frantically with their superiors to find answers needed for crucial decisions. Meanwhile workers stand idle, waiting for directions from above. While it is not possible in general to have all needed information available from one's own subordinates, concurrent systems require this locality property to reduce interference from too much communication.

The design of computers and of algorithms has yet to show the ingenuity reflected in human organizations. This failure is not for want of cleverness in designers, but rather because the technologies used to implement computers are much less flexible than the human beings used to implement corporations. VLSI offers more flexibility than earlier technologies because memory and processing structures can be implemented with the same technology, in close proximity.

Communication costs in computers

The archetypal computer consists of a single "processor" (the CPU or "central processing unit"), connected to a large, homogeneous memory (Figure 1). The processor fetches an instruction from
memory, decodes it, executes it, and repeats the cycle. Many instructions will cause additional references to memory in order to fetch operands or to store results. The performance of such a computer depends critically on the speed with which memory can be accessed.

A very simple argument can be developed to determine the speed of the memory. If a memory of \( M \) bits is implemented on a single chip in a two-dimensional array, wires approximately \( M^{1/2} \) long are required to transmit data between a memory cell and the processor. (We are concerned with relative units of length and time, because we intend only to compare different designs, not to determine absolute execution speeds.) The time required for data transmission is proportional to this length: the longer the wire, the greater the distance the signal must propagate and the greater the wire's capacitance, slowing propagation. In addition to slowing the memory, long wires also consume a great deal of chip space and require substantial power to drive. In present implementations of large computers, performance is further decreased by the several levels of packaging required to provide a memory of significant size: chip, printed-circuit board, backplane. The wiring on chips and printed-circuit boards grows as \( M^{1/2} \), but backplane wiring grows linearly with memory size.

The organization shown in Figure 1 is also rather wasteful of resources: most of the memory and memory wiring is idle most of the time. For a typical large memory, \( M \) might be \( 32 \times 10^6 \), but only a 16- or 32-bit word will be delivered to the processor with each memory reference. If the memory is organized as an array of \( 10^6 \) bits for each bit in the word, only 2 of the 2000 wires needed to address the array are used in a given reference (1000 select wires running horizontally, and 1000 select wires running vertically). Vast areas of memory thus lie idle because the amount of information extracted on a single reference is small compared to the size of the entire memory.

The costs of communication are exorbitant in today's computers. Most of the expense, time, and energy required to compute are consumed by the communication of data over large distances.

Memory Locality

Computer designers have recognized the difficulty of communicating with a very large memory, and have taken steps to utilize the memory more effectively. The result is a memory hierarchy, outlined in Figure 2. The processor communicates with a series of memories, whose size increases and speed decreases as they become farther from the processor. The closest memory (\( M_F \)) provides high-speed "registers" or "accumulators" that are used very frequently, usually to
contain intermediate results of arithmetic calculations. Next comes "cache" memory ($M_c$),
designed to hold data and instructions that are referenced frequently. The "primary" memory
($M_p$) is similar to the large memory of several million bits illustrated in Figure 1. Finally, a
"secondary" memory ($M_s$) of some sort is provided, usually implemented with disks.

The average time required to reference a memory element will depend on which piece of the
memory hierarchy holds the desired element. The intent is that fast, small memories be
referenced more frequently than the slow, large ones. This desire is reflected in the design of the
instruction set of the computer: referencing "registers" is usually encouraged by the structure of
the instruction set; referencing primary memory (or cache) is supported by the instruction set, but
perhaps in less flexible ways than for register access; finally, accessing a disk is not directly
supported by instruction sets at all, but requires complicated "I/O control."

It is instructive to formulate a crude model to estimate the performance of the memory hierarchy.
We need to assign representative values to the frequency with which each memory is accessed,
and to the size of each memory:

\[
\begin{align*}
  f_r &\sim 0.6 & \text{Frequency of access to registers ($M_r$)} \\
  f_c &\sim 0.38 & \text{Frequency of access to cache ($M_c$)} \\
  f_p &\sim 0.02 & \text{Frequency of access to primary memory ($M_p$)} \\
  f_s &\sim 0.000005 & \text{Frequency of access to secondary memory ($M_s$)} \\
  M_r &\sim 16 \\
  M_c &\sim 10^3 \\
  M_p &\sim 10^5 \\
  M_s &\sim 10^{10}
\end{align*}
\]

Using our model of memory access time, the time required to access memory on the average is

\[f_r M_r + f_c M_c + f_p M_p + 100 f_s M_s\]

measured in arbitrary units. (The factor of 100 arises because disk access times are substantially worse than our memory wiring model indicates.)

It is instructive to note the relative contributions of the separate memories: 2.4, 12, 6, 50, for a
total of 70. The cost of access to the slowest memory, the disk, is the most important
contribution to the average.

The memory hierarchy is an improvement over the homogeneous memory of Figure 1. The time
to reference a single memory of size $10^5$ is 320 units. The time to reference a three-level
hierarchy of about the same size \((M_P, M_C, M_P, \text{ with frequencies shown above})\) is a mere 20 units.

The effectiveness of the memory hierarchy depends on *locality* of the memory references. Cache algorithms copy large chunks (8-32 words) of primary memory into the cache, hoping that additional memory references will occur in the neighborhood of the first reference. A similar hope is attached to transfers from secondary memory. If an application arises in which most of the memory references do *not* go to the fast register memory, the memory hierarchy will perform poorly.

Locality can also be viewed as a function of size. If a program and its data can reside in primary memory for the duration of execution, and do not require secondary memory, the average memory access time will drop from 70 to 20. If the program is small enough to fit in the small cache memory, access time will drop further to 14.

**Concurrency in computers**

Not content with the increases in speed due to a memory hierarchy, computer designers have also sought to increase the concurrency in computer designs. A number of different approaches have been tried (see reference [3]); we shall illustrate pipeline structure and multiprocessor structures.

**Pipelined processors**

Pipelined processors are patterned after the production line found in manufacturing: a portion of the processing is performed by each of several processors, and then handed to the next processor in the line. Starting from Figure 1, the designer reasons that two processors could function concurrently, each assigned to half the original memory (Figure 3); a communication path is provided so that the first processor can transmit results to the second.

The two-processor pipeline more than doubles the processing power available. If we neglect the cost of inter-processor communication, the time required to execute an instruction is \((1/2)^{1/2}(M/2)^{1/2}\), about one third the time required by the uniprocessor in Figure 1. The improvement comes from two effects: doubling the number of processors doubles the speed, but reducing the memory size also increases speed.

A special case of pipelining is illustrated by *instruction-fetch overlap* in computers. One processor is responsible for fetching an instruction from memory; it then passes on to the second processor...
information required to execute the instruction; the second processor actually performs the execution. In chapter 6, we saw this technique applied in OM: while one microinstruction is being executed, the controller is fetching the next microinstruction. Execution overlap allows the execution itself to be pipelined (see reference [6] for more pipelining structures).

Pipelined structures are perhaps most effective in special-purpose applications that can utilize a large number of processors. Signal-processing is a particularly good example: a signal is sampled digitally to generate a stream of signal data. This data is pipelined through processors to perform corrections, correlations, frequency analysis, etc. Section 2 of this chapter illustrates the application of pipelines to matrix arithmetic of various sorts.

Unfortunately, it is not always possible to cast problems in a framework suited to execution on pipelined computers. If the workload is not divided evenly among the processors, some will stand idle, reducing the effective speed increase. But it is the rigid communication discipline that most severely restricts the application of pipelines.

Multiprocessors

Another important class of concurrent computers are multiprocessors. Unlike the pipeline, these structures provide switching structures that allow each processor to communicate with each other processor. The hope is that those algorithms not suited to pipelines because of their communication requirements can be executed on multiprocessors.

Figure 4 shows a dual-processor configuration, again adapted from Figure 1. Each processor communicates primarily with a memory half the size of the original. In addition, a common "bus" is provided to allow each processor to reference the other's memory.

Two problems with the dual-processor arrangement are immediately apparent. First, if each processor references memories at random, the two will interfere often, and vitiate some of the speed gain. Second, can we assure that the sequential program suited to the uniprocessor architecture of Figure 1 can be adapted to the dual-processor configuration? Putting aside for the moment the problems of programming a multiprocessor, we shall examine its performance.

We shall construct a crude model of the time required to execute an instruction on the dual processor. Assume that each processor references its own memory with probability \((1/p)\), and the other's with probability \(p\). Further, assume that the useful duty cycle of each processor is \(d\). If
both processors can be productively employed at all times, \( d \) will be 1. However, if the two processors must occasionally wait for each other, i.e., must "synchronize," \( d \) may fall below 1. We can identify three cases:

1. \( P_a \) references \( M_a \) and \( P_b \) references \( M_b \); probability is \((1-\beta)^2\).
2. \( P_a \) and \( P_b \) both reference \( M_a \) (or, equivalently \( M_b \)); probabilities sum to \( 2\beta(1-\beta) \).
3. \( P_a \) references \( M_b \) and \( P_b \) references \( M_a \); probability is \( \beta^2 \).

We also need to model the time required to complete each of the three cases. A processor references its own memory, of size \( M/2 \), in time \((M/2)^{1/2}\). When a reference is made to a neighbor's memory, we assume the time for communication on the bus and referencing the memory sum to \#M, as if it were addressing the entire memory as one array. The costs for the three cases then become:

1. \((M/2)^{1/2}\)
2. \((M/2)^{1/2} + M^{1/2}\)
3. \(M^{1/2} + M^{1/2}\)

From these estimates we calculate the expected instruction execution time, remembering that \( 2d \) processors are available:

\[
\text{time} = M^{1/2} \left( \frac{1}{d} \right) \left( 2^{1/2}/4 + \beta \cdot \beta^2/2 \right)
\]

This expression is plotted in Figure 5, assuming \( d=1 \).

The simple model of a dual-processor configuration is suggestive of behavior we can expect from multi-processor systems that require global communication. We observe that if \( \beta = 0 \), execution speed is more than twice that of the uniprocessor illustrated in Figure 1. Just as in the pipeline, doubling the number of processors contributes a factor of two, but additional speed is achieved because each processor addresses a smaller memory.

The model also illustrates the importance of locality in the use each processor makes of its memory. If \( \beta \) is allowed to grow too large, the factor of two contributed by two processors is erased by interference between the processors when accessing the common memory.

Perhaps the most important parameter is \( d \), which is determined by our ability to adapt algorithms to multi-processor configurations. Some applications seem to decompose nicely for execution on concurrent hardware, and some offer difficulties. In human organizations we have become resigned to *always* attacking large problems in a concurrent way. We will, no doubt, have to do the same with computer programs.
Summary

The schemes we have illustrated that reduce communication costs and try to exploit concurrency can be combined in various ways in computer structures. The table below summarizes the speedup effect that these techniques offer, as derived from our crude models ($n$ denotes the number of processors used):

<table>
<thead>
<tr>
<th>Technique</th>
<th>Typical speedup factor</th>
</tr>
</thead>
<tbody>
<tr>
<td>Memory hierarchy</td>
<td>10</td>
</tr>
<tr>
<td>Pipelining</td>
<td></td>
</tr>
<tr>
<td>instruction overlap</td>
<td>2</td>
</tr>
<tr>
<td>special-purpose</td>
<td>$n$</td>
</tr>
<tr>
<td>Multithread processors</td>
<td>$&lt; n$</td>
</tr>
</tbody>
</table>

The processor-memory structures and algorithms presented in the remainder of this chapter all attempt to have as many processors as can be kept productive simultaneously and to locate them as close as possible to the data they require. These are the considerations exhibited by our simple models of memory hierarchies, pipelines and multiprocessors. The examples presented here by no means exhaust the topic of concurrent computation; the interested reader will find literatures on computer architecture [2,3], parallel processor and processing [4,5,6,7], performance evaluation [3], and algorithm design [8,9,10,11,14].
2. Algorithms for VLSI Processor Arrays

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2.1. Introduction

"And the smooth stream in smoother numbers flows"

--Alexander Pope

We are interested in high-performance parallel algorithms that can be implemented directly on low-cost hardware devices. By performance, we are not referring to the traditional operation counts that characterize classical analyses of algorithms, but rather, the throughput obtainable when a special purpose peripheral device is attached to a general purpose host computer. This implies that time spent in I/O, control, and data movement as well as arithmetics must all be considered. The cost of the device must be measured in how well it can be implemented using LSI technology and must be sensitive to what the technology can do cheaply, and what is expensive.

LSI technology has made one thing clear. Simple and regular interconnections lead to cheap implementations and high densities, and high density implies both high performance and low overhead for support components. The two-dimensional array structure consisting of mesh-connected processors enjoys this desirable property. Therefore, we are interested in designing parallel algorithms which have simple and regular data flows so that they can be executed efficiently on such processor arrays. We are also interested in using pipelining as a general method for implementing these algorithms in hardware. By pipelining, processing may proceed concurrently with input and output, and consequently overall execution time is minimized. Pipelining plus multiprocessing at each stage of a pipeline should lead to the best-possible performance. In this section, we demonstrate simple and regular multiprocessor networks that are capable of pipelining some important matrix computations with optimal speed-up.
Most of the results reported here are based on a paper by H. T. Kung and C. E. Leiserson, which is to be presented at the Symposium on Sparse Matrix Computations and Their Applications in Knoxville, TN, November 2-3, 1978.

2.2. The Basic Components and Structures

The single operation common to the problems considered in this section is the so-called inner product step, \( C = C + A \times B \). We postulate a processor which has three registers \( R_A, R_B, \) and \( R_C \). Each register has two connections, one for input and one for output. Fig. 2.2.1 shows two types of geometries for this processor. Type (a) geometry will be used for matrix-vector multiplication and solution of triangular linear systems (Sections 2.3 and 2.6), whereas type (b) geometry will be used for matrix multiplication and LU-decomposition (Sections 2.4 and 2.5).

![Diagram of processor geometries](image)

Fig. 2.2.1. Geometries for the inner product step processor.

The processor is capable of performing the inner product step. We shall define a basic time unit in terms of this processor. At time \( t \), the processor shifts its inputs into \( R_A, R_B, \) and \( R_C, \) and computes \( R_C = R_C + R_A \times R_B \). At time \( t+1 \), the new value of \( R_C \) together with the input values for \( R_A \) and \( R_B \) are available as outputs. All outputs are latched and the logic is
clocked so that when one processor is connected to another, the changing output of one during
the time interval between $t$ and $t+1$ will not interfere with the input to the other during this
time. This is not the only processing element we shall make use of, but it will be the work
horse. These special processors will be specified later when they are used.

The basic network organization we shall adopt for internal communication is the mesh-connected
processor scheme. (See Fig. 2.2.2.) All connections from a processor are to neighboring
processors. The most widely known system based on this organization is the ILLIAC IV. If
diagonal connections are added in one direction only, we shall call the resulting scheme
hexagonally mesh-connected or hex-connected for short. We shall demonstrate that linearly
connected and hex-connected processors are natural for matrix problems.

![Diagram](attachment:diagram.png)

**Fig. 2.2.2.** Examples of mesh-connected processors.

When an input path to a processor lies on an edge of the device, we shall sometimes designate
it as an external input connection from the host memory. Alternatively, we may let the input
have a fixed value such as zero. An output data path will either go to the host memory or be
ignored.
2.3. Matrix-Vector Multiplication

We consider the problem of multiplying a matrix with a vector. Let $A = (a_{ij})$ be an $n \times n$ band matrix with band width $w = p + q - 1$, and $x = (x_1, ..., x_n)^T$, $y = (y_1, ..., y_n)^T$ be $n$-vectors such that $Ax = y$. (See Fig. 2.3.1 for the case when $p = 2$ and $q = 3$.)

![Matrix-vector multiplication diagram](image)

Fig. 2.3.1. The matrix-vector multiplication when the matrix is a band matrix with $p = 2$ and $q = 3$.

Suppose $A$ and $x$ are given. The following algorithm computes the product $y = Ax$ by pipelining the computation through $w$ linearly connected processors. Before giving the code for each processor, we illustrate the algorithm for the band matrix-vector multiplication problem in Fig. 2.3.1. For this case the linearly connected network has four processors. See Fig. 2.3.2.
Fig. 2.3.2. The linearly connected network for the matrix-vector multiplication problem shown in Fig. 2.3.1.

The general scheme of our pipelining algorithm can be viewed as follows. The \( y_i \), which are initially zero, keep moving to the left while the \( x_i \) are moving to the right and the \( a_{ij} \) are moving down. It turns out that each \( y_i \) is able to accumulate all its terms, namely, \( a_{i,i-2}x_{i-2}, a_{i,i-1}x_{i-1}, a_{i,i}x_i, \) and \( a_{i,i+1}x_{i+1} \), before it leaves the network. Fig. 2.3.3 illustrates the first seven steps of the algorithm. Note that when \( y_1 \) and \( y_2 \) are output they have the correct values. Observe also that at any given time alternating processors are idle. (Indeed, it is possible to use \( w/2 \) processors in the network for a general band matrix with band width \( w \). We did not do so for the sake of clarity.)

We now specify the algorithm more precisely. Assume that the processors are numbered by integers 1, 2, \ldots, \( w \) from the left end processor to the right end processor. Each processor has three registers, \( R_A \), \( R_x \) and \( R_y \), which will hold entries in \( A \), \( x \) and \( y \), respectively. Initially, all registers contain zeros. Each step of the algorithm consists of the following operations, but for odd numbered time steps only odd numbered processors are activated and for even numbered time steps only even numbered processors are activated.
1. *Shift.*

- $R_A$ gets a new element in the band of matrix $A$.

- $R_x$ gets the contents of register $R_x$ from the left neighboring node. (The $R_x$ in processor 1 gets a new component of $x$.)

- $R_y$ gets the contents of register $R_y$ from the right neighboring node. (Processor 1 outputs its $R_y$ contents and the $R_y$ in processor $w$ gets zero.)

2. *Multiply and Add.*

$$R_y = R_y + R_A \times R_x.$$  

Using the processor postulated in section 2.2, we note that the three shift operations in step 1 can be done simultaneously, and that each step of the algorithm takes a unit of time. Suppose the bandwidth of $A$ is $w = p + q - 1$. It is readily seen that after $w$ units of time the components of the product $y = Ax$ start shifting out from the left end processor at the rate of one output every two units of time. Therefore, using our network all the $n$ components of $y$ can be computed in $2n + w$ time units, as compared to the $O(wn)$ time needed for the sequential algorithm on a single processor.
Fig. 2.3.3. The first seven steps of the matrix-vector multiplication algorithm.
2.4. Matrix Multiplication

This section considers the problem of multiplying two matrices. Let $A$ and $B$ be $n \times n$ band matrices of bandwidth $w_1$ and $w_2$, respectively. We show that a network of $w_1 w_2$ hex-connected processors can compute the product $C = AB$ in $3n + \min(w_1, w_2)$ units of time. The algorithm uses the same principle as the one in Section 2.3. We illustrate the general scheme by considering the matrix multiplication problem depicted in Fig. 2.4.1.

The diamond shaped interconnection network for this case is shown in Fig. 2.4.2, where processors are hex-connected and data flows are indicated by arrows. The nonzero elements in $A$, $B$ and $C$ move through the network in three directions, as indicated in the figure. Initially, the $c_{ij}$ are all zeros. One can easily see that with the type (b) inner product processors described in Section 2.2, each $c_{ij}$ is able to accumulate all its terms before it leaves the network.

Suppose that Fig. 2.4.2 describes the configuration at time 1. Then, for example, $c_{11}$ gets $a_{11}b_{11}$ at time 2 and $a_{12}b_{21}$ at time 3, while $c_{21}$ gets $a_{21}b_{11}$ at time 3 and $a_{22}b_{21}$ at time 4. (Note that approximately only one third of processors in the network are active at a given time. Indeed, it is possible to use about $(w_1 w_2)/3$ processors in the network for multiplying two band matrices with band widths $w_1$ and $w_2$.)
Fig. 2.4.1. Matrix multiplication.

Fig. 2.4.2. The network for the matrix multiplication $C = A \times B$ shown in Fig. 2.4.1.
2.5. The LU-Decomposition of a Matrix

The LU-decomposition of a given a matrix $A$ is the problem of computing lower and upper triangular matrices $L$ and $U$ such that $A = LU$. (Cf. Fig. 2.5.1.)

![LU-Decomposition Diagram]

Fig. 2.5.1. The LU-decomposition of a matrix.

Once the $L$ and $U$ factors are known it would be relatively easy to solve a linear system $Ax = b$ or to invert $A$. In the following we describe a parallel LU-decomposition algorithm using a hex-connected network.

We assume that $A$ is either a symmetric positive-definite or irreducible diagonally dominant matrix. It is well-known that under this assumption the $L$ and $U$ matrices can be obtained by Gaussian elimination without pivoting. We show the rather surprising fact that Gaussian elimination enjoys the same data flow as matrix multiplication and that all the processors except one perform the same inner product step. In fact, the same matrix multiplication network in Section 2.4 can be used to compute $L$ and $U$ matrices, provided that the processor at the top now computes minus the reciprocal of an input and the orientation of the other boundary processors is properly altered. More precisely, at the special processor at the top, the data from the south processor is passed unchanged to the north, minus its reciprocal is computed and sent to the southwest processor, and a numerical value "1" is sent to the southeast processor. The processors on the left hand "upper" side are rotated 120 degrees clockwise and always receive "0" from their northwest external input connections. Similarly, the processors on the right hand "upper" side are rotated 120 degrees counterclockwise and always receive "0" from their northeast external input connections. (Of course, it is not necessary to actually input "0" for
these processors; we did so for the sake of uniformity.

Suppose that $L = (l_{ij})$ and $U = (u_{ij})$. Then Gaussian elimination computes the entries in $L$ and $U$ using the following procedure:

$$l_{ij} = m_{ij} \text{ for } i > j, \ 1 \text{ for } i = j \text{ and } 0 \text{ for } i < j,$$

$$u_{ij} = a_{ij}^{(i)} \text{ for } i \leq j \text{ and } 0 \text{ for } i > j,$$

$$a_{ij}^{(0)} = a_{ij},$$

$$m_{ik} = a_{ik}^{(k)} / a_{kk}^{(k)},$$

$$a_{ij}^{(k+1)} = a_{ij}^{(k)} + m_{ik} a_{kj}^{(k)}.$$

To illustrate our results, we consider a band matrix $A$ with $p = 4$ and $q = 4$. When elements in the band of matrix $A$ are fed into the lower edge of the hex-connected network as shown in Fig. 2.5.2, the elements of $L$ and $U$ are output from the upper edge. Fig. 2.5.3 shows an enlargement of the configuration after eight steps of the algorithm have been executed. The flow of data on the network is indicated by arrows in Fig. 2.5.3. The hexagons denote the standard processors which perform the inner product step just like the corresponding processors in the matrix multiplication network (cf. Fig. 2.4.2). The processor at the top denoted by a circle performs the reciprocal and negation operations. As in the matrix multiplication algorithm, each processor only operates once every three time steps. We will not give a formal correctness proof for the algorithm here. But for understanding the algorithm the reader is advised to view the LU-decomposition as the inverse problem of multiplying a lower triangular matrix with 1's on the diagonal to an upper triangular matrix. Then the algorithm of this section can simply be regarded as one which undoes the matrix multiplication algorithm of Section 2.5. Having realized this, one should be able understand also why the two algorithms use the same network and enjoy the same data flow pattern. The idea of using the same network for both the forward and backward problems seems to be general. It will be used again in Section 2.6.
Fig. 2.5.2. The hex-connected network for pipelining the LU-decomposition of a band matrix with $p = 4$ and $q = 4$. 
Fig. 2.5.3. LU-decomposition after the first eight steps.

It is readily seen that if matrix A is nxn, then using the network shown in Fig. 2.5.2 the L and U matrices can be computed in 3n+4 units of time. In general, if A is an nxn band matrix with band width w = p+q-1, then with a network of no more than pq hex-connected processors, the LU-decomposition of A can be done in 3n+min(p,q) units of time. (It is possible to reduce the number of required processors to about pq/3.) In particular if A is an nxn dense matrix, then n² hex-connected processors can compute the L and U matrices in 4n units of time, including I/O time.
2.6. Triangular Linear Systems

Suppose that we want to solve a linear system $Ax = b$. Then after having done with the LU-decomposition of $A$ (e.g., by methods described in Section 2.5), we still have to solve two triangular linear systems $Ly = b$ and $Ux = y$. This section concerns itself with the solution of triangular linear systems.

Let $A = (a_{ij})$ be a nonsingular $n \times n$ band lower triangular matrix with band width $w=q$. Suppose that $A$ and an $n$-vector $b=(b_1,\ldots,b_n)^T$ are given. The problem is to compute $x=(x_1,\ldots,x_n)^T$ such that $Ax=b$. (See Fig. 2.6.1 for the case when $q=4$.)

![Triangular Matrix](image)

Fig. 2.6.1. The band (lower) triangular linear system with $q=4$.

We show that this problem can be solved by the algorithm and network almost identical to those used for band matrix-vector multiplication in Section 2.3. (Note that the linear system problem can be regarded as the inverse of the matrix-vector multiplication problem.) We illustrate our result by considering the linear system problem in Fig. 2.6.1. For this case, the network and the general scheme of the algorithm are described in Fig. 2.6.2.
The $y_i$, which are initially zero, keep moving to the left while the $x_i$, $a_{ij}$ and $b_i$ are moving in the network, as indicated in Fig. 2.6.2. The left end processor is special in that it performs $x_i^- (b_i - y_i)/a_{ii}$. Each $y_i$ accumulates inner product terms in the rest of the processors as it moves to the left. At the time $y_i$ reaches the left end processor it has the value $a_{i1}x_1 + a_{i2}x_2 + ... + a_{ini}x_{i-1}$, and, consequently, the $x_i$ computed by $x_i^- (b_i - y_i)/a_{ii}$ at the processor will have the correct value. Fig. 2.6.3 demonstrates the the first ten steps of the algorithm.
<table>
<thead>
<tr>
<th>Step Number</th>
<th>Configuration</th>
<th>Comments</th>
</tr>
</thead>
<tbody>
<tr>
<td>0</td>
<td></td>
<td>$y_1$ enters processor 4.</td>
</tr>
<tr>
<td>1</td>
<td></td>
<td>$y_1$ moves left one position.</td>
</tr>
<tr>
<td>2</td>
<td></td>
<td>$y_2$ enters processor 4.</td>
</tr>
<tr>
<td>3</td>
<td></td>
<td>$x_1 = (b_1 - y_1)/a_{11}$, $(x_1 = b_1/a_{11}$, since $y_1 = 0$.)</td>
</tr>
<tr>
<td>4</td>
<td></td>
<td>$y_2 = a_{21}x_1$.</td>
</tr>
<tr>
<td>5</td>
<td></td>
<td>$x_2 = (b_2 - y_2)/a_{22}$.</td>
</tr>
<tr>
<td>6</td>
<td></td>
<td>$y_3 = a_{31}x_1$.</td>
</tr>
<tr>
<td>7</td>
<td></td>
<td>$x_3 = (b_3 - y_3)/a_{33}$.</td>
</tr>
<tr>
<td>8</td>
<td></td>
<td>$y_4 = a_{41}x_1 + a_{42}x_2 + a_{43}x_3$.</td>
</tr>
<tr>
<td>9</td>
<td></td>
<td>$x_4$ is output.</td>
</tr>
</tbody>
</table>

Fig. 2.6.3. Solving a lower band triangular system with $q = 4$.

One can check that the computed $x_1$, $x_2$, $x_3$ and $x_4$ all have correct values. With this network we can solve an $n \times n$ band triangular linear system with band width $w = q$ in $2n + q$ units of time.
2.7. Applications and Comments

2.7.1 Variants of the Algorithms

Rather than the basic algorithms presented above it is their variants that will be used mostly in practice. No attempt is given here for listing all the possible variants; it is important that the reader understands the basic principles of the algorithms so that he can construct appropriate variants for his specific problems.

We first note that although most of our illustrations are done for band matrices all the algorithms work for the regular nxn dense matrix. In this case the band width of the matrix is w = 2n - 1. If the band width of a matrix is so large that a corresponding algorithm requires more processors than a given network provides, then one should decompose the matrix and solve each subproblem on the network.

One can often reduce the number of processors required by an algorithm if the matrix is known to be sparse. For example, the matrices derived from differential equations by using finite differences or finite elements approximations are usually “sparse band matrices.” These are band matrices whose nonzero entries appear only in a few of those lines in the band which are parallel to the diagonal. In this case by introducing proper delays to each processor for shifting its data to its neighbors, the number of processors required by the algorithms in Sections 2.3 and 2.6 can be reduced to the number of those diagonal lines which contain nonzero entries. This variant is useful for performing iterative methods involving sparse band matrices.

It is possible to use our algorithms and networks to solve some nonnumerical problems when appropriate interpretations are given to the addition (+) and multiplication (x) operations. For example, some pattern matching problems can be viewed as matrix problems with comparison and Boolean operations.

2.7.2 Convolution and Discrete Fourier Transform

There are a number of important problems which can be formulated as matrix-vector multiplication problems and thus can be solved rapidly by the algorithm in Section 2.3. The problems of computing convolutions and discrete Fourier transforms are such examples. If a matrix has the property that the entries on any line parallel to the diagonal are all the same, then
the matrix is a Toeplitz matrix. The convolution problem is simply the matrix-vector multiplication where the matrix is a triangular Toeplitz matrix (cf. Fig. 2.7.1).

\[
\begin{bmatrix}
  a_1 \\
  a_2 & a_1 \\
  a_3 & a_2 & a_1 \\
  a_4 & a_3 & a_2 & a_1 \\
  \vdots & \vdots & \vdots & \vdots & \ddots
\end{bmatrix} 
\begin{bmatrix}
  x_1 \\
  x_2 \\
  x_3 \\
  x_4 \\
  x_5 \\
  \vdots \\
  \vdots \\
  \vdots \\
  \vdots
\end{bmatrix} = 
\begin{bmatrix}
  b_1 \\
  b_2 \\
  b_3 \\
  b_4 \\
  b_5 \\
  \vdots \\
  \vdots \\
  \vdots \\
  \vdots
\end{bmatrix}
\]

Fig. 2.7.1. The convolution of vectors \(a\) and \(x\).

On the other hand the \(n\)-point discrete Fourier transform is the matrix-vector multiplication, where the \((i,j)\) entry of the matrix is \(\omega^{(i-1)(j-1)}\) and \(\omega\) is a primitive \(n\)th root of unity (cf. Fig. 2.7.2).

\[
\begin{bmatrix}
  1 & 1 & 1 & 1 \\
  1 & \omega & \omega^2 & \omega^3 \\
  1 & \omega^2 & \omega^4 & \omega^6 \\
  1 & \omega^3 & \omega^6 & \omega^9 \\
  \vdots & \vdots & \vdots & \vdots & \ddots
\end{bmatrix} 
\begin{bmatrix}
  x_1 \\
  x_2 \\
  x_3 \\
  x_4 \\
  x_5 \\
  \vdots \\
  \vdots \\
  \vdots \\
  \vdots
\end{bmatrix} = 
\begin{bmatrix}
  b_1 \\
  b_2 \\
  b_3 \\
  b_4 \\
  b_5 \\
  \vdots \\
  \vdots \\
  \vdots \\
  \vdots
\end{bmatrix}
\]

Fig. 2.7.2. The discrete Fourier transform of vector \(x\).
Therefore using a linearly connected network of size $O(n)$ both the convolution of two $n$-vectors and the $n$-point discrete Fourier transform can be computed in $O(n)$ units of time, rather than $O(n \log n)$ as required by the sequential FFT algorithm. Moreover, note that for the convolution problem each processor has to receive an entry of the matrix only once, and this entry can be shipped to the processor through horizontal connections and stay in the processor during the rest of the computation. For the discrete Fourier transform problem each processor can in fact generate on-the-fly the powers of $\omega$ it requires. As a result, for these two problems it is not necessary for each processor in the network to have the external input connection on the top of the processor, as depicted in Fig. 2.3.2.

In the following we describe how the powers of $\omega$ can be generated on-the-fly during the process of computing an $n$-point discrete Fourier transform. The requirement is that if a processor is $i$ units apart from the middle processor then at time $i + 2j$ the processor must have the value of $\omega^{i^2 + ij}$ for all $i, j$. This requirement can be fulfilled by using the algorithm below. We assume that each processor has one additional register $R_t$. All processors except the middle one perform the following operations in each step, but for odd (respectively, even) numbered time steps only processors which are odd (even) units apart from the middle processor are activated. For all processors except the middle one the contents of both $R_A$ and $R_t$ are initially "0".

1. *Shift.* If the processor is in the left (respectively, right) hand side of the middle processor then

   - $R_A$ gets the contents of register $R_A$ from the right (respectively, left) neighboring processor.
   
   - $R_t$ gets the contents of register $R_t$ from the right (respectively, left) neighboring processor.

2. *Multiply.*

   $$R_A \leftarrow R_A \times R_t.$$
The middle processor is special; it performs the following operations at every even numbered time step. For this processor the contents of both $R_A$ and $R_t$ are initially "$1".

1. $R_A \leftarrow R_A \times R_t^2 \times \omega$.
2. $R_t \leftarrow R_t \times \omega$.

2.7.3. The Common Memory Access Pattern

Note that all the algorithms given in this section retrieve and store elements of the matrix in the same order. (See Fig. 2.3.2, 2.4.2, 2.5.2, and 2.6.2.) Therefore, we recommend that matrices be always arranged in memory according to this particular ordering so that they can be accessed efficiently by any of the algorithms.

2.7.4. The Pivoting Problem

In Section 2.5 we assume that the matrix $A$ has the property that there is no need of using pivoting when Gaussian elimination is applied to $A$. What should one do if $A$ does not have this nice property? (Note that Gaussian elimination becomes very inefficient on mesh-connected processors if pivoting is necessary.) This question motivated us to consider Givens' transformation for triangularizing a matrix, which is known to be a numerically stable method. It turns out that, like Gaussian elimination without pivoting, the orthogonal factorization based on Givens' transformation can be implemented naturally on mesh-connected processors, although a pipelining implementation appears to be more complex.

2.8. Concluding Remarks

Research in interconnection networks and algorithms has been traditionally motivated by large scale array computers such as ILLIAC IV (see, for example, Kuck[5] and Stone [3]). The results presented in this section were, however, motivated by the advance in integrated circuit technology, though they are certainly applicable to parallel array computers. We have shown that many basic matrix computations can be done very efficiently by special purpose multiprocessors, which may be built cheaply using the current technology. The common feature of our algorithms is that their data flows are very simple and regular, and they are pipeline algorithms. We have discovered that some data flow patterns and interconnection schemes are fundamental for matrix computations. For example, the two-way flow on the linearly connected network is common to
both matrix-vector multiplication and solution of triangular linear systems (Sections 2.3 and 2.6), and the three-way flow on the hexagonally mesh-connected network is common to both matrix multiplication and LU-decomposition (Sections 2.4 and 2.5). A practical implication of this fact is that one device may be used for solving many different problems. Moreover, we note that almost all the processors needed in any of these devices are the inner product processor postulated in Section 2.2. A careful design for this processor is desirable since it is the work horse for all the devices presented.

For the important problem of solving a dense system of \( n \) linear equations in \( O(n) \) time on \( n \times n \) mesh-connected processors, we have improved upon the recent results of Kant and Kimura [13]. The basis of their results is an theorem on determinants which was known to J. Sylvester in 1851. Their algorithm requires that the matrix be "strongly nonsingular" in the sense that every square submatrix is nonsingular. It is sufficient for our algorithms in Section 2.5 that the matrix be symmetric positive-definite or irreducible diagonally dominant.

We end this section by noting that processor communication will likely continue to dominate the cost of parallel algorithms and systems. Communication paths inherently take more space and energy than processing elements. We regard the problem of minimizing communication costs as fundamental. We hope the results of this section have demonstrated that the communication problem in parallel algorithms is not only tractable but also interesting. We expect that a large number of algorithms having small communication costs will be discovered in the future.
3. Hierarchically Organized Machines

We know that human organizations use hierarchical structure to extract the greatest possible benefit from the daily activities of tens of thousands of individuals. We know that complex systems can be constructed by subdividing them into less complex systems, which are again subdivided, as many times as necessary until the resulting systems are simple enough to construct easily. We have seen that the organization of real estate on the silicon surface dictates a hierarchical communication system for any device which must support global communication. Such hierarchical communication exists in conventional computers only in a limited way. Are there new machine structures which communicate hierarchically, which support systems that consist of an arbitrary hierarchy of subsystems, and which can coordinate the activities of any number of submachines?

Binary Trees

Consider any number of processors physically arranged as a binary tree. Each processor has two subprocessors which it can control. These subprocessors, in turn, have two sub-subprocessors, and so on. A possible layout of such a binary processor tree is shown in figure 8. At the lowest level a small array of ordinary memory cells, labeled $M_0$, is accessed by the lowest level processors, labeled $P_0$. The combination of one lowest level processor with its associated memory is the element of computing power. These units are grouped together in pairs and accessed by the next level processor, labeled $P_1$. Two $P_1$'s with their associated lower level units are grouped together and accessed by the next level higher processor, labeled $P_2$. This arrangement is repeated recursively until an entire silicon chip is covered by the processor memory hierarchy. The rate at which information can be transferred within a processor is independent of the level of the processor. As the wires within a processor get longer, the drivers must become proportionately larger to drive them. The highest level processor which communicates off the silicon chip to the outside world has large drivers and hence is able to drive off chip without suffering a severe performance penalty. Such a machine can thus be extended to a large number of individual chips and still maintain the full speed of the individual processors within it.

A conventional computer is a special case of this organization, consisting of a memory cell and a bottom-level processor. Also, there is another way to map a conventional computer onto a binary tree of processors. View the highest level processor as a cpu and load all subprocessors with programs that merely decode requests for the memory below them. Loaded with these programs,
the structure between the two extreme levels becomes a memory decoder tree between a conventional CPU and its memory.

More importantly, this binary tree structure is a completely general, concurrent processing engine and can be used for problems decomposed in an arbitrary hierarchical way. If a problem requires more than two subprocessors at any level, a subtree of physical processors can be operated as one logical processor, matching the problem's structure. Algorithms for constructing logical processors of any size are given in the next section. The tree has inherent in it the ability for all processors to compute concurrently and hence has a vastly larger potential computing power than a conventional machine using a similar amount of silicon real estate.

Since the number of processors decreases exponentially with the level, the total bandwidth available, whether processing or communication, decreases exponentially with the level. Half of the total bandwidth of the system is concentrated at level 0, one quarter at level 1, one eighth at level 2, etc. A particular computation is well matched to such a processor if its bandwidth requirements are concentrated at the lowest levels. If an algorithm requires more communication at any level than the structure provides, it will not be able to take advantage of all the processing power of the structure. An extreme example of this sort is the von Neumann machine where all computation occurs at the highest level processor and the lower level processors are used only once at a time as an ordinary memory. Such a machine requires equal bandwidth at each level of the hierarchy and is an exponential waste of the resources of the machine.

It is also clear that such a structure is testable if a single processor is testable. Each supervisor merely loads a test program into its two subordinates and exercises them. Once it has established that both work correctly, it loads each with the program it just used to test them. A tree of $N$ levels can thus be tested in $N$ times the time necessary to test one processor.

It is difficult to predict how any radically different machine structure will perform in a real computing environment. Ideally, one should implement a number of complete systems, spanning a large range of user requirements, in order to gain experience with the strengths and weaknesses of any given scheme. Failing that, we can at least map certain algorithms onto our machine in the hope that they will shed light on its capabilities and its problems. Several such mappings are presented in the next section. We plan to develop others and we hope our readers will contribute still more for subsequent editions of this text.
Algorithms for the Tree Machine

Section contributed by Sally Browning, Caltech

A. A Word About Notation

The notation chosen to describe the processor tree and the algorithms that run in each node of the tree must emphasize that the number of different flavors of processors is small (usually one). That is, a few templates describe them all.

Secondly, we want to emphasize locality. The processor tree is interesting because each node is a powerful computing engine that can work independent from its neighbors. Our notation must be one that encourages self-sufficient modules.

We will use a modified version of the SIMULA syntax. The CLASS concept of SIMULA provides us a means of describing a template that will be instantiated as the nodes of the tree. We can designate individual procedures and data elements as either local to this node or visible to the outside world. And we can use recursion to indicate flow of execution through the tree.

Most importantly, though, SIMULA’s CLASS construct is designed for expressing and enforcing locality. SIMULA is an object-oriented language, and, as such, encourages the programmer to think in terms of objects doing operations to themselves. The knowledge of the representation and meaning resides in the class, not in some omnipotent overlord. This is exactly the notation we need to describe the nodes of our tree.

Because we are describing highly concurrent algorithms we need to get around the sequential nature of SIMULA statements. We expand the meaning of the semicolon symbol. In conventional SIMULA, semicolon is used to terminate a statement. We use semicolon to make a statement about the execution as well. Read semicolon as “At this point, all statements in progress must be terminated before advancing to the next statement.” Linefeed will be used to indicate syntactic end of the statement. In other words, linefeeds are used to separate statements; semicolons are used to separate groups of statements which can execute concurrently.
B. A Word About Branching Ratios

While the physical structure of our tree restricts each processor to two descendants, we can impose a logical structure that allows an arbitrary branching ratio. Each logical processor consists of several physical processors, enough to provide the desired number of offspring. A logical node with $N$ children is built from $N-1$ physical nodes and is $\log N$ levels deep. Figure 1 shows some sample logical processors.

We can describe the process of mapping our logical structure onto the physical tree in SIMULA. We define two CLASSes: a node and a processor. A node represents the physical entity. It has exactly two descendants. A processor will refer to the logical entity, with an arbitrary number of children.

In the SIMULA definitions, $N$ represents the number of descendants desired. As we build the logical node, we attempt to keep it balanced. That is, all available physical nodes on a given level of the tree will be used before a new level is added. Nodes on a given level are added to the logical processor from left to right, as in Figure 1. Note that CLASS Processor is a refinement of CLASS Node that knows how to choose one of $N$ descendants.

```
CLASS Node(n): INTEGER n;
BEGIN
  REF(Node)left, right;

  // logic code to build logical node:
  IF n > 2 THEN left := NEW Node((n + 1)/2);
  IF n > 3 THEN right := NEW Node(n/2);
END of CLASS Node;

Node CLASS Processor;
BEGIN

  REF(Processor) PROCEDURE Son(s): INTEGER s;
BEGIN
  p := IF s <= (n + 1)/2 THEN left ELSE right;
  WHILE NOT (p IN Processor) DO
    p := IF s <= (p.n + 1)/2 THEN p.left ELSE p.right;
  Son := p;
END of PROCEDURE Son;

END of CLASS Processor:
```
Figure 1. Logical Nodes (solid color) with Two to Seven Descendants

Figure 2. Systematic Generation of Subgraphs in a Graph of 4 Nodes
II. Algorithms with Polynomial Complexity

One of the traditional approaches to solving a problem that is too large or too complex when considered as a whole is to break the problem recursively into pieces that are manageable. The point is to apply as many concurrent processors to the problem as possible in order to reduce execution time. We will look at two algorithms that use this approach, sorting and matrix multiplication. While both of these problems are solved nicely on cellular arrays, it is instructive to map them onto a machine with different communication properties.

A. Sorting in linear time.

We use a binary tree with depth $\log N$ to sort $N$ numbers. The sort is accomplished as a byproduct of loading the numbers into memory and then reading them out again. The numbers themselves are never in sorted order internally, but come out of the tree in the desired order.

This algorithm is an implementation of heap sorting, one of the well known techniques used in sequential machines [14]. It is a particularly interesting example because it illustrates a fundamental issue in concurrency. It is well known that sorting on a sequential machine can be done with $O(N \log N)$ comparisons. Heap sorting requires $O(N^2)$ comparisons, and has been considered inferior for that reason. However, it has been shown on very fundamental grounds that if communication is restricted to nearest neighbors, at least $N^2$ comparisons are required [17]. The apparent advantage of the $O(N \log N)$ algorithms comes as a direct result of longer communication paths. It is also clear that no scheme will be able to produce an ordered set of numbers until all numbers to be sorted are loaded into the machine. For this reason, the best we can expect is to use $N$ processors for $O(N)$ cycles.

The algorithm that runs in each processor has a procedure for loading the tree called Fillup and a procedure invoked during the output cycle called Passup. Fillup keeps the largest number seen to date, and passes the smaller one to the right or left child, keeping the tree balanced by alternating sides. Passup returns this processor’s current number and refills itself with the larger of the numbers stored in its descendants. This action is pipelined so that the largest number is always available in the root.
Below is a SIMULA description of the algorithm running in each processor. The variable number holds the number stored in this processor. The boolean symbol empty reflects the validity of that number. The boolean identifier balanced is used to keep the tree balanced as it is loaded.

CLASS processor;
BEGIN
  INTEGER number;
  BOOLEAN balanced, empty;
  REF(processor) left, right;

  PROCEDURE fillup(candidate): INTEGER candidate;
  BEGIN
    IF empty THEN
      BEGIN
        number := candidate
        empty := FALSE;
      END
    ELSE
      BEGIN
        IF candidate > number THEN  !swap;
        BEGIN
          t := candidate;
          candidate := number;
          number := t;
        END;
        IF balanced
        THEN left.fillup(candidate)
        ELSE right.fillup(candidate);
        balanced := NOT balanced;
      END;
  END of procedure fillup;

  INTEGER PROCEDURE passupnumber;
  BEGIN
    passupnumber := number;
    IF left = NONE AND right = NONE THEN empty := TRUE  !its a leaf;
    ELSE
      IF left.empty THEN
        BEGIN
          IF right.empty THEN empty := TRUE  !left & right subtrees empty;
          ELSE number := right.passupnumber;  !fill from right son;
        END
      ELSE
        BEGIN
          IF right.empty THEN number := left.passupnumber  !fill from left son;
          ELSE number := IF left.number > right.number
          THEN left.passupnumber ELSE right.passupnumber;
          !take the larger of the two;
        END;
    END
  END of procedure passupnumber;

  !init code:
  empty := FALSE;
  balanced := TRUE;
  !left and right set;

END of class processor;
B. Matrix Multiplication.

Suppose we have two NxN matrices to multiply together. By using the divide and conquer approach, we can break the problem down until we have \( N^3 \) problems that multiply two numbers (1x1 matrices) together. We reassemble the matrix on the way back up the tree.

We use the following rule to subdivide the problem:

Let A, B, and C be NxN matrices such that \( AB = C \). We subdivide all three into four \( N/2 \times N/2 \) submatrices, e.g., \( A_{11}, A_{12}, A_{21}, A_{22} \).

Then \( C = ( A_{11} B_{1j} + A_{i2} B_{2j} ) \), \( i,j = 1,2 \).

We will consider matrices of size \( N = 2^M \) without loss of generality. A tree to multiply two matrices of size \( 2^M \) will have \( M \) levels of processors that add two matrices together, \( M \) levels that split and assemble the matrix, and one level (the leaf nodes) that multiply two numbers together. Thus the tree is \( 2M+1 \) logical levels deep.

Each adder node has two descendants, and each split/assemble node has four descendants. Thus the physical structure will use two levels to simulate the 4-way branching, and the tree will, in fact, be \( 3M \) levels deep. That is, the tree is \( 3 \log N \) levels deep and therefore has \( N^3 \) leaf nodes. Thus a total of \( 2N^3 \cdot 1 \) processors are used in the computation.

Let us look at the communication requirements between nodes of the tree. The root node must be prepared to store the entire matrix. The adder nodes in level one (the root is level 0) will each deal with a quarter of the original matrix, as will the split/assemble nodes in level three. The further down the tree you go, the smaller the matrix the node must store and transfer.

However, note that each of the \( N^2 \) elements must travel the entire length of the tree and back again during the execution of the algorithm. While communication requirements are low at the leaves, they are extremely high (roughly \( N^2 \) numbers to receive, and \( (N/2)^2 \) to pass down to each descendant) at the root.

In the algorithm given below, the add operation takes \( N^2 \) time. By splitting this up among parallel processors by row (\( N \) of them) or element (\( N^2 \) of them) we can make this operation linear or constant in time. However, the problem is still limited by the split/assemble process that requires each element to travel the height of the tree. That is, the best time performance we can achieve with this algorithm is \( N^2 \log N \).
We now present a SIMULA representation of a matrix and use it in the algorithm that follows.
The algorithm uses two kinds of processors, the adders and the split/assemble nodes. Each matrix is divided into submatrices as follows:

\[
\begin{array}{cc}
1 & 2 \\
3 & 4 \\
\end{array}
\]

CLASS matrix(n); INTEGER n;
BEGIN

INTEGER ARRAY val[1:n,1:n];

REF(matrix) PROCEDURE quarter(select); INTEGER select;
BEGIN REF(matrix) =q:
    INTEGER i,j,k,l;
    qa:=NEW matrix(n//2);
    i:=j:=1;
    IF select = 2 THEN j:=n//2+1
    ELSE IF select = 3 THEN i:=n//2+1
    ELSE IF select = 4 THEN i:=j:=n//2+1;
    FOR k:=1 STEP 1 UNTIL qa.n DO
        FOR l:=1 STEP 1 UNTIL qa.n DO
            qa.val[k,l]:=val[i+k*1,j+l-l-1];
    quarter:=qa;
END of procedure quarter:

REF(matrix) PROCEDURE compose(a,b,c,d); REF(matrix) =a,b,c,d;
BEGIN INTEGER i,j;
    FOR i:=1 STEP 1 UNTIL a.n DO
        FOR j:=1 STEP 1 UNTIL a.n DO
            val[i,j]:=a.val[i,j];
        FOR i:=1 STEP 1 UNTIL b.n DO
            FOR j:=1 STEP 1 UNTIL b.n DO
                val[i,j+a//2]:=b.val[i,j];
        FOR i:=1 STEP 1 UNTIL c.n DO
            FOR j:=1 STEP 1 UNTIL c.n DO
                val[i+n//2,j]:=c.val[i,j];
        FOR i:=1 STEP 1 UNTIL d.n DO
            FOR j:=1 STEP 1 UNTIL d.n DO
                val[i+n//2+j+n//2]:=d.val[i,j];
    compose:=THIS matrix;
END of procedure compose;
END of class matrix;
CLASS processor(size);
BEGIN

REF(matrix)mat;
REF(processor)one,two,three,four;

REF(matrix) PROCEDURE multiply(a,b); REF(matrix)a,b;
BEGIN REF(matrix)c;
    c::NEW matrix(a.n);
    IF c.n = 1 THEN
        c.val[1,1] := a.val[1,1]*b.val[1,1]
    ELSE
        c.compose(one.multiply(a.quarter(1),b.quarter(1),a.quarter(2),b.quarter(2)),
                   two.multiply(a.quarter(3),b.quarter(3),a.quarter(4),b.quarter(4)));
    END;
    multiply:=c;
END of procedure multiply;

REF(matrix) PROCEDURE mult&add(a,b,c,d); REF(matrix)a,b,c,d;
BEGIN REF(matrix)c1,c2; INTEGER i,j;
    c1::one.multiply(a,b);
    c2::two.multiply(c,d);
    FOR i := 1 STEP 1 UNTIL c1.n DO
        FOR j := 1 STEP 1 UNTIL c2.n DO
            c1.val[i,j] := c1.val[i,j] + c2.val[i,j];
        END OF procedure mult&add;
END of class processor;
III. Solutions to Nonpolynomial Problems.

Complexity theory [9,10] has established a context within which it is possible to make certain statements concerning the inherent complexity of computations. These statements are universally couched in the terminology of sequential machines. There is, however, a class of problems for which the possibility of large scale concurrency has been addressed.

Consider a computation in which there are N conceptual steps. At each step, q alternative branches may be taken. Such a computation may be viewed as a tree with $q^N$ possible outcomes. If at each step there is enough information available to decide which branch to take, a sequential machine will be able to complete the computation in KN cycles where K is the average number of cycles spent at each step. The dependence of the number of machine cycles upon the number of conceptual steps is thus linear. The problem is said to be linear in N or of order N, written O(N).

In many computations, not enough information has been generated by previous steps to determine which branch to take. Later steps will generate this information, but we cannot execute the later steps until after the earlier steps! In such cases, the sequential machine must simply try one branch at random. If it concludes after executing subsequent steps that the particular branch taken was wrong, it must backtrack to the original point, and try another route.

In a wild flight of fancy, we might become frustrated with this behavior and wish we had a machine which was so smart that it could tell if it was on the right path, even if there was no possibility of choosing such a path with the information at hand. It would make an arbitrary choice at each branch—and always be right! Such a machine cannot, of course, be built with real logic operating with real programs. However, we can imagine such a machine in much the same way we imagine a space ship traveling faster than the speed of light. Machines of this sort are called nondeterministic, since there is no way this behavior can be specified on rational grounds.

Returning to our problem, it is clear that a sequential nondeterministic machine could solve the problem in O(N) cycles. Problems which can be solved by such an imaginary nondeterministic machine in a number of cycles which is bounded by some fixed power of N are said to be Non deterministic-Polynomial abbreviated NP [19,9].

It is quite clear that the behavior of a nondeterministic machine can be simulated by a set of concurrent deterministic machines. Each machine can simply follow a separate path through the
tree. At the end, there will be $q^N$ processors, representing each possible outcome of the computation. Although different problems will have different branching ratios ($q$) and different depths ($N$), all can be mapped onto the tree machine using techniques described earlier.

It has been shown that there is a class of problems of this sort where there are no shortcuts. Working one path through to the end gives no clue concerning the outcome of another path. Such problems are, in some sense, maximally difficult. They are called *NP-complete* problems.

A great deal of lore has developed concerning NP-complete problems. It has been shown that, in some sense they are all "equivalent" [18]. Suppose machine Y can solve a single kind of NP-complete problem. The equivalence property states that there is an algorithm which will run on an ordinary sequential machine in a polynomial number of cycles that transforms a description of any NP-complete problem into a description of a problem solvable by Y. If Y can solve its NP-complete problems in polynomial time, then it can be used to solve any NP-complete problem in polynomial time. If Y requires exponential time, any NP-complete problem will also require exponential time.

The methods we use to describe trees of different branching ratios to a binary tree machine are very similar to the methods used to map an NP-complete problem onto a machine that solves another. When a tree with branching ratio greater than 2 is mapped onto a binary tree, the depth of the tree increases. Mapping a tree with branching ratio less than 2 will decrease the depth. In a similar fashion, the algorithm that transforms NP-complete problems may increase the number of alternative branches ($q$) and decrease the number of conceptual steps ($N$) or vice-versa. Thus the mappings that establish the equivalence class of NP-complete problems are exactly like the mappings from trees of one branching ratio to another.

The theory that establishes the NP-complete equivalence class offers direct guidance in mapping such problems onto a highly concurrent structure. Because we can solve any one problem in our concurrent tree machine, and because we know a mapping from an arbitrary NP-complete problem into this one, we can solve the arbitrary problem.

The traditional approach to solving the class of problems that grow exponentially has been to recognize space or processing power as a limited resource. The problems have exponential time complexity because the solutions proceed sequentially. As VLSI becomes a reality, however, it is interesting to treat processors as an unlimited resource and look at the time complexity of these
problems when they take advantage of concurrency. We emphasize, however, that while the time complexity is significantly reduced, problems require an exponential number of processors. If you solve a problem of reasonable size, you will use an enormous number of processors. In a later section, an example is worked for an NP-complete problem that grows as $N^N$. The problem uses a graph of 4 nodes, and our concurrent solution requires 95 processors. A graph of 10 nodes could use as many as $2^{10^{10}}$ processors!

We will examine two NP-complete problems. The clique problem has time complexity of $O(2^N)$ when the possible cliques are considered sequentially. The color cost problem is $O(N^N)$. By taking advantage of the parallel consideration of possible solutions, using $O(2^N)$ and $O(N^N)$ processors respectively, we will present solutions to these two problems that take polynomial, in fact $O(N^2)$, time.
A. The Clique Problem.

A clique is complete subgraph. That is, given an undirected graph G, a clique C contained in G is a graph such that for all nodes n,m in C, there is an edge (n,m). Finding the largest clique in an arbitrary graph is an NP-complete problem.

Given a graph G with N nodes, numbered from 1 to N, we will consider each node sequentially and generate potential cliques. Ignoring the edges for a moment, a collection of M nodes leads to $2^M - 1$ potential cliques. This, interestingly enough, is the number of nodes in a binary tree of depth M. We will use this fact to generate the cliques in our graph incrementally.

Each level in the tree represents the addition of another node to be considered. Each processor at a given level will spawn two descendants. The left child will consider the subgraph consisting of the new node and all but the last node of the parent subgraph. The right child’s subgraph will add the new node to the complete parent subgraph. In this manner, we generate all possible subgraphs for a graph of N nodes. Figure 2 is an example for N=4.

If each node stores an edge list, the tree can be pruned of subgraphs that are not cliques. The number of processors required is reduced, but the worst case behavior is identical. At most $2^{N-1}$ processors are required to solve the problem for a graph of size N. Our solution, which uses pruning, requires $O(N^2)$ time.

Each processor stores the edge list as a boolean matrix called edge, an integer size that holds the size n (number of nodes) of the clique this processor represents. An array called clique contains the numbers of the nodes that form the clique.

When a processor is activated, by a call to the procedure FindClique, it will already have a clique assigned to it. FindClique’s purpose is to generate cliques for its descendent nodes. It does this according to the method described above. That is, if the subgraph that contains the new node and all of the nodes in clique except the last one is a clique, it will be assigned to the left child. Likewise, if the addition of node to clique yields a complete subgraph, the right child will represent it. If either of the subgraphs is not complete, the descendant will not be generated.

The tree of all cliques is generated iteratively by considering each node of the graph in turn. In the main program given below, p is a reference to the root processor. Each processor in the tree will pass up the largest clique among its children. Thus the root returns the size of the largest clique known to date.
CLASS processor;
BEGIN
  REF(processor).left,right;
  BOOLEAN ARRAY edge[1:n,1:n];
  INTEGER ARRAY clique[1:n];
  INTEGER size;

  BOOLEAN PROCEDURE IsClique;
  BEGIN INTEGER i,j;
    IsClique:=TRUE;
    FOR i:=1 TO size DO
      FOR j:=1 TO size DO
        IF NOT edge[i,j] THEN IsClique:=FALSE;
      END of procedure IsClique;
    END;

    REF(processor) PROCEDURE FindClique(node);
    BEGIN INTEGER r,REF(processor).l,r;
      l:=r:=THIS processor;
      IF size=0 THEN
        BEGIN
          clique[i]:=node
          size:=1
          FindClique:=THIS processor;
        END;
      ELSE IF left.size=0 THEN
        BEGIN
          FOR i:=1 TO size-1 DO left.clique[i]:=clique[i]
          left.clique[size]:=node
          left.size:=size;
          IF NOT left.IsClique THEN left.size:=0;
        END;
      ELSE l:=left.FindClique(node);
      IF right.size=0 THEN
        BEGIN
          FOR i:=1 TO size DO right.clique[i]:=clique[i]
          right.clique[size+1]:=node
          right.size:=size+1;
          IF NOT right.IsClique THEN right.size:=0;
        END;
      ELSE r:=right.FindClique(node);
      IF l.size > r.size
        BEGIN
          FindClique:=l
        END;
      ELSE IF r.size > size
        THEN FindClique:=r;
      ELSE FindClique:=THIS processor;
    END;
  END;
END of procedure FindClique:

size:=0;
!left and right set up correctly;
!read in edge list;
END of class processor;

!main program to start it all up:
BEGIN REF(processor).largest:=INTEGER i
  FOR i:=1 TO n DO largest:=p.findclique(node);
END of main:
Figure 3. Sample Graph for Clique Problem

Figure 4. Tree built to find Cliques in Graph of Figure 3
Figure 3 gives a sample graph of six nodes. Figure 4 shows the processor tree that is built and used to find the cliques in the graph. The tree has height 6, and the largest cliques have 4 nodes. Each processor in the tree represents a clique in the graph.

B. The Color Cost Problem.

This NP-complete problem is an adaptation of the K-colorability problem. Given an undirected graph $G$ of $N$ nodes and a set of $N$ colors, each with an associated cost, we want to find a minimum cost coloring of the graph such that no nodes sharing an edge are the same color.

There are $N^N$ possible colorings of the graph. Evaluating them sequentially produces a solution in time $O(N^N)$. We present a parallel algorithm that requires $O(N^2)$ time and $O(N^N)$ processors.

In this problem we will make use of the ability to simulate arbitrary branching ratios on our binary tree. We will discuss the problem in terms of logical nodes with up to $N$ descendants. An earlier part of this section describes the method of mapping logical structures onto the physical one.

As in the clique problem, each level in the processor tree represents the consideration of another node. That is, level one shows possible colors for the first node, level two colors the second node based on the choices made for at level one, and so on. We will describe the generation of the potential colorings.

Each node has an edge list called edge and a list of costs indexed by color number called colorcosts. There is an array called coloring that reflects the color choices for preceding nodes, and a boolean array called colors that is used to generate the possible colorings for this node.

The algorithm, given in procedure color begins by assuming that all colors yield valid colorings. The array coloring is used to eliminate those colors that have been used to color nodes that share an edge with this node. This reduced set of colors, all of which are legal colorings, is used to spawn descendants, one for each coloring of this node.

When the tree is $N$ levels deep all the legal colorings have been generated. The leaf nodes calculate a cost for the coloring they represent, and each parent node takes as its cost the least cost among its children. Thus the minimum cost coloring is stored at the root.
Here is the algorithm that will run in each processor.

CLASs processor;
BEGIN

    BOOLEAN ARRAY edge[1:n,1:n],colors[1:n];
    INTEGER ARRAY coloring[1:n],colorcost[1:n];
    INTEGER cost;

    PROCEDURE color(node): INTEGER node;
    BEGIN INTEGER i;
        IF node > n THEN
            BEGIN
                cost := 0;
                FOR i := 1 TO node-1 DO cost := cost + colorcost[coloring[i]];
            END
        ELSE
            BEGIN
                FOR i := 1 TO node-1 DO IF edge[i,node] THEN
                    colors[coloring[i]] := FALSE;
                FOR i := 1 TO n DO
                    IF colors[i] THEN
                        BEGIN
                            son(i).coloring[node] := i;
                            son(i).color(node+1);
                        END
                    ELSE son(i) := NONE;
                cost := maxcost;
                FOR i := 1 TO n DO
                    IF (IF son(i) = NONE THEN FALSE ELSE cost > son(i).cost) THEN cost := son(i).cost;
            END;
    END of procedure color;

    REF(Processor) PROCEDURE Son(s): INTEGER s;
    BEGIN REF(node):=
        p:=IF s <= (n+1)/2 THEN left ELSE right;
        WHILE NOT (p IN Processor) DO
            p:=IF s <= (p.n+1)/2 THEN p.left ELSE p.right;
        Son:=p;
    END of PROCEDURE Son;

END of class processor;

Let us work a small example. We will use the graph and color set given in Figure 5. Figure 6 shows the colorings and costs arrived at by the algorithm. Each level of the tree represents a node of the tree. That is, if the root is level 0, the first node is colored in level 1, and level 4 represents potential colorings for the fourth node. Besides representing a part of a coloring, each node also contains the minimum cost coloring found among its descendent colorings.
Figure 5. Sample Graph and Color Table for Color-Cost Problem

<table>
<thead>
<tr>
<th>Color</th>
<th>Cost</th>
</tr>
</thead>
<tbody>
<tr>
<td>blue</td>
<td>0</td>
</tr>
<tr>
<td>green</td>
<td>1</td>
</tr>
<tr>
<td>red</td>
<td>2</td>
</tr>
<tr>
<td>black</td>
<td>3</td>
</tr>
</tbody>
</table>

Figure 6. Color-Cost Tree for the Example of Figure 5
We see that there are two equivalent colorings that yield the minimum cost of 3. Coloring nodes (1,2,3,4) either (green,blue,red,blue) or (red,blue,green,blue) gives us a coloring with minimum cost.

V. Conclusions

The tree of processors we have described is a general computing structure. Each node in the tree is a processor with general computing capability. It is not designed with a specific problem or class of problems in mind.

The most dramatic results are achieved when the machine is applied to a problem that can take advantage of the concurrency the tree of processors provide. We have presented solutions to four problems that, in varying degrees, have this characteristic.

The four examples we have presented in this section can be summarized by citing the execution time and the number of processors required. Note that the total chip area of a tree machine is related to the number of processors.

<table>
<thead>
<tr>
<th>Problem</th>
<th>Time</th>
<th>Processors</th>
</tr>
</thead>
<tbody>
<tr>
<td>Sorting</td>
<td>$N$</td>
<td>$N$</td>
</tr>
<tr>
<td>Matrix multiplication</td>
<td>$N^2 \log N$</td>
<td>$2N^3-1$</td>
</tr>
<tr>
<td>Clique</td>
<td>$N^2$</td>
<td>$2N-1$</td>
</tr>
<tr>
<td>Color cost</td>
<td>$N^2$</td>
<td>$N^N$</td>
</tr>
</tbody>
</table>

If an algorithm exhibits exponential growth, as do the clique and color-cost problems, the lower bound on time complexity is $N$. A tree with an exponential number of leaves will be $O(N)$ deep. Again, our solutions do not realize this lower bound. The loading of the edge matrix is an $O(N^2)$ operation. Additionally, each node of the graph is considered in turn, and causes the traversal of a tree of depth up to $N$. This too is of $O(N^2)$ in time. Are there better algorithms that can achieve the lower bound complexity?

Because we are used to designing machines for a sequential environment, we do not yet understand the effect that concurrency will have on the conceptualization of problem solutions. An open question is to characterize those problems that can benefit from the concurrency provided by our tree of processors. Are the communication paths of the tree adequate for this
class of problems? Can we design algorithms with the traditional programming notations, or does
their sequential nature hide the concurrency? Can NP-complete problems be solved in O(N) time
with an unlimited number of processors? What can be said about the concurrency of NP-
complete problems in general? These are just a few of the interesting questions that arise from
the study of a concurrent environment.
4. Highly Concurrent Structures with Global Communication

(Adapted from a paper by Carver Mead and Martin Rem [16])

This section presents an analysis of the constraints placed by physical laws on a VLSI system in which information must be communicated from any location to any other. The spectacular performance of cellular arrays on algorithms which map into regular structures leads us to ask if there is any possibility that a highly concurrent structure can be built which will act as a general purpose computational engine. Such generality must include the ability to transmit information over distances as large as required by the computation. Before describing such a machine we will analyze in detail the requirements that global communication places on the design of any computing structure.

There has previously been no adequate theoretical basis for optimizing the overall organization of systems implemented in the VLSI technology. Conventional complexity theory is inadequate because its measure of cost is the number of steps taken by a sequential machine to complete the computation. No account is taken of the size of the machine (and hence the time required for each step). Possible concurrency is ignored, thereby ruling out the most important potential contribution of the silicon technology. Traditional switching theory is also inadequate. While it provides a beautiful formalism for describing elementary logic functions, its optimization methods concern themselves with logical operations rather than communication requirements. Even in today’s integrated circuits, the wires required for communicating information across the chip account for most of the area. Driving these wires accounts for most of the time delay and energy dissipation. In very large scale integrated systems, the situation becomes even more extreme. In this section, we describe a method by which the conceptual organization of a large chip can be analyzed, and a lower bound placed on its size, cycle time, and energy dissipation, before a detailed design is undertaken. The results of this analysis suggest rather general guidelines for the organization of all large integrated systems.
Metrics of Space, Time, and Energy

Physical Properties

Devices used to construct monolithic silicon integrated circuits are universally of the charge-controlled type. A charge $Q$ placed on the control electrode (gate, base, etc.) results in a current $I = Q/\tau$ flowing through the device. The transit time $\tau$ is the time required for charge carriers to move through the active region of the device.

All times in an integrated system can be formulated as simple multiples of $\tau$. For one transistor to drive another identical to it, a charge $Q$ must flow through its active region, requiring time $\tau$. If the capacitance $C_L$ of the load driven is $K$ times the gate capacitance $C_g$ of the driving transistor, a time $K\tau = C_L/C_g \tau$ is required. Likewise, the elementary energy associated with the signal charge $Q$ on the gate capacitance $C_g$ is $E_0 = C_g V^2/2$. A load capacitance $K C_g$ requires an energy $K E_0$. Since wires have a minimum width, their capacitance is directly proportional to their length. Thus the energy required to transmit a signal from one point on the chip to another is proportional to the distance separating the two points. As the unit of length we employ the minimum spacing of two conducting paths. For the unit of time we choose the time it takes a minimum size transistor to charge a wire of unit length plus another transistor like itself. One unit of time is thus slightly larger than the transit time of a transistor.

Advantages of Hierarchical Structures

We are considering large integrated systems in which it is necessary to communicate information throughout the entire system. As an example, consider a bit of information stored on the gate of a minimum size transistor in a random-access memory which must be communicated to the memory bus of a CPU. Since there are many words of data in the memory, there are many possible sources for each wire in the memory bus. Figure 1 illustrates two possible approaches to organizing such a bus. In the first approach, a transistor associated with each bit drives the bus wire directly. If the bus wire has a capacitance $C_w$, the time required to drive the bus wire is $t = \tau C_w/C_g$. In a typical computer memory $C_w$ is many orders of magnitude larger than $C_g$, and the delay introduced by such a scheme is very long. Since $C_w$ is proportional to the length of the wire, it is also proportional to $S$, the number of driver transistors connected to the wire and $b$, the spacing between transistors. Assuming most of the capacitance $C_w$ is due to the wire itself;
Fig 1a
A bus driven directly by memory cells

Fig 1b
A bus driver tree
Fig. 2
Delay of a hierarchical structure as a function of alpha
\[ t = b_0 \tau S \quad (1) \]

A second scheme is shown in Figure 1b. Here each transistor drives a wire only long enough to reach its neighbor. Each such wire is connected to the gate of a transistor twice as large as the transistor driving it. The arrangement is repeated upward until the top level where all sources have a path to the bus. In this scheme the delay in driving the lowest level wire is approximately \( 2\tau b_0 \). The delay introduced by the wires at each level is the same, since each driver transistor is twice as large as those driving it. Hence the delay in driving the bus line is \( 2\tau N b_0 \) where \( N \) is the number of levels in the structure. Since there are \( S = 2^N \) transistors at the lowest level, the delay may be written:

\[ t = 2\tau b_0 \log_2 S \quad (2) \]

Comparing (2) and (1), we see that for large \( S \) the delay has been made much shorter by using a hierarchical structure.

**A Cost Criterion**

A hierarchy such as that shown in Figure 1b may use any integral number, \( \alpha \), of transistors driving each wire. We refer \( \alpha \) as the branching ratio of the driver hierarchy. The driver transistors will in general be \( \alpha \) times the size of those driving them. The delay for such a structure is \[ t = \alpha \tau b_0 \log_\alpha S = b_0 \tau \alpha / \log_\alpha, \] dependent upon the branching ratio of the hierarchy. This delay is plotted in Figure 2, normalized to its minimum value which is attained at \( \alpha = e \).

While dramatic improvements in the performance of integrated structures can be achieved by a hierarchical organization, a penalty is always paid in the area required for wires. In the simple case shown, a bus requiring one wire when driven directly requires \( \log_\alpha S \) wires when organized as a hierarchy. For this reason it is not possible to optimize a design without a cost function involving both area and time. In this paper we will use the area-time product as an example of such a cost function. Other cost functions may be more appropriate under some circumstances. For the above simple example, the cost function is area \(*\) time = \( b_0 \tau (\log S)^2 \alpha / (\log \alpha)^2 \). This cost is minimized for \( \alpha = e^2 \approx 7.4 \).

**Hierarchical computing systems**

The analysis given above suggests a very general structure for computing systems. Lowest level
cells are grouped together into modules in such a way that \( \alpha \) cells drive their outputs onto an output wire. Each output wire is connected to a driver transistor which is \( \alpha \) times as large as those driving the wire. Modules are grouped in such a way that \( \alpha \) of those modules drivers are connected to an inter-module communication wire. This wire in turn is connected to a driver transistor \( \alpha^2 \) times as large as the lowest level transistors. This process is continued until the appropriate size system has been realized. Notice that the area of the driver transistor for each wire in such a structure is proportional to the area of the wire. For this reason, we compute only the area of the wires. The drivers somewhat enlarge the unit of wire area, but do not change the functional form of the solutions.

**Random-Access Memory - an example**

In this section we discuss the design of a large of a random-access memory (RAM) of \( S \) bits. We will apply a rigid structural discipline to our design, and compute the cost and performance of the resulting memory.

**Organization of the RAM**

We organize the RAM in a hierarchical fashion. The elements of level 0 are the bits themselves, each bit consisting of two crossing wires: a select wire and a data wire. When the select wire is asserted, it puts its contents on the data wire. We group \( \alpha^2 \) bits into an \( \alpha \times \alpha \) square to form a module of level 1. If the width of an element (a bit) is \( b_0 \) the elements have to drive wires of length \( \alpha b_0 \). A module on level 1 consists of an array of horizontal select and vertical data wires, constituting the \( \alpha^2 \) bits of level 0, and some additional logic and wires at the side. We group again \( \alpha^2 \) of these modules into a square to form a module of level 2, etc. Figure 3 shows three levels of the hierarchy for \( \alpha = 4 \).

To study the memory in more detail we look at a module of level 1 (Figure 4). We describe how one extracts one of its \( \alpha^{2i} \) bits. In order to select one bit of storage \( 2i \log_\alpha \) address wires are required. We run \( i \log_\alpha \) of them, called the row address wires, vertically along the side of the module and the other \( i \log_\alpha \), the column address wires, horizontally. Its \( \alpha^2 \) submodules are organized into \( \alpha \) rows of \( \alpha \) submodules each. When the select wire of the module is asserted \( \log_\alpha \) of the row address wires are used, by the decoder, to select one of the \( \alpha \) rows of submodules; the select wire running through that row is asserted. The other \((i-1)\log_\alpha \) row address wires are run horizontally into each of the \( \alpha \) rows of submodules, where they serve as
Fig. 3

Three Levels of a Memory Hierarchy with $\alpha = 4$
Fig. 4
A RAM module of level i (i > 0)
column address wires for the submodules. Of the \( \log \alpha \) column address wires \((i-1)\log \alpha\) are run vertically into each of the \( \alpha \) columns of submodules, where they serve as row addresses. The other \( \log \alpha \) address wires are used by the multiplexer to select one of the \( \alpha \) data wires coming out of the columns of submodules. The signal on the selected data wire is driven onto the data wire of the module itself.

If we wish to have a memory of \( S \) bits with \( N + 1 \) levels (level 0 through \( N \)) we choose \( N = \frac{\log S}{2\log \alpha} \), or \( S = \alpha^{2N} \). A hierarchical structure results which contains \( S \) bits from which we can extract one bit at a time. If we want the word length to be \( W \) we employ \( W \) of these structures in parallel: to select one word we select one bit in each of the \( W \) hierarchies.

**Area of the RAM**

Figure 4 allows us to compute the size of a RAM. Let \( L_i \) denote the width of a module of level \( i \), then we have the following recurrence relation:

\[
L_0 = b_0 \\
L_i = i\log \alpha + 1 + \log \alpha + \alpha L_{i-1}
\]

The solution to the above relation is

\[
L_i = \alpha^i b_0 + (\alpha^i - 1)/(\alpha - 1) + (2\alpha^{i+1} - \alpha^i \alpha)/((\alpha - 1)^2 - i + 1/\alpha - 1)! \log \alpha
\]

Rather than the width itself we are interested in the width per bit. In one direction, horizontal or vertical, module \( i \) has \( \alpha^i \) bits: we therefore compute \( L_i/\alpha^i \).

\[
L_i/\alpha^i = b_0 + 1/(\alpha - 1) + 2\alpha - 1/(\alpha - 1)^2 \log \alpha - 1/(\alpha - 1) \alpha^i [(\alpha/\alpha - 1 + 1 + i) \log \alpha + 1]
\] (3)

An interesting property of the width per bit, as expressed by (3), is that its limit for \( i \to \infty \) is finite.

\[
\lim_{i \to \infty} L_i/\alpha^i = b_0 + 1/\alpha - 1 + 2\alpha - 1/(\alpha - 1)^2 \log \alpha
\] (4)

This means that the width per bit \( L_i/\alpha^i \) is bounded from above by (4) independent of the number of levels of a RAM. Expression (3) converges in an exponential fashion towards its limit: for small values of \( i \) (3) is already very close to (4). We, therefore, use (4) as the width per bit for a RAM; its square is then the area per bit. By dividing the area per bit by the bit area \( b_0 \).
we obtain the total area per bit area for a RAM. Figure 5 shows this quotient as a function of $\alpha$ for four different values of $b_0$. It gives the overhead factor in the area that is due to the wires. A memory chip will be larger by this factor than the area of its level 0 cells alone. For a memory of 64K bits with $N=2$, $\alpha$ should be 16. Expression (4) is then equal to $b_0 + 0.6$. This shows that in 2-level 64K dynamic MOS memories, for which $b_0$ lies between 1 and 2, roughly half of the area will be occupied by wires.

One may wonder why we have not discussed the area that is consumed by the wires for power and ground. The reason is that these wires can be thought of as increasing only the width $b_0$ of each bit; they do this by an amount that is roughly independent of $\alpha$, as is shown in the following analysis.

For simplicity we assume that the wires for power and ground run in opposite directions, say parallel to the data and select wires. We compute how much one of them contributes to the width of a module $i$. The width of a power or ground wire is proportional to the number of bits served by it. Let the width of the highest level be $u$, given $S$ and the design of the lowest level memory cell this parameter is easy to compute. The width of the wire in a module on level $i$ is proportional to the current it must supply and is hence $u \alpha^i / \alpha^2$. In one direction, horizontal or vertical, there are $\alpha^N / \alpha^i$ such modules. The total contribution of all modules on level $i$ is thus $u \alpha^i / \alpha^N$. Taking the sum of this expression for $i=0,1,...,N$ yields $u / \alpha^N \alpha^{N+1}/\alpha-1 \approx u \alpha/\alpha-1$. There are $u/S$ bits in one direction, the increase of the bit width, due to power and ground, is therefore

$$u/S \alpha/\alpha-1,$$

which is roughly equal to $u/S$.

We are interested in the optimal choice of $\alpha$, but to make that choice we will have to look at the access time, which also depends on $\alpha$.

*Access time of the RAM*

Each element of level 0 drives a wire of length $\alpha b_0$ to reach the periphery of its module on level 1; this takes time $\alpha b_0$. Each module on level 1 drives in the same amount of time a wire that is $\alpha$ times longer to reach the periphery of its module on level 2, etc. With $N$ being the level of the highest module, the time required to extract one bit of storage adds up to $\alpha b_0 N$. We use
Fig. 5  Total area per bit of a RAM as a function of alpha
Fig. 6 Area-time product of a RAM as a function of alpha
this figure as the access time. For a RAM of \( S \) bits the access time is then \( a b_0 \log S/2 \log \alpha \).

**Cost of the RAM**

We take the product of the area and the access time as the cost function of the RAM. A RAM of \( S \) words of \( \log S \) bits each has the following area-time product.

\[
(b_0 + 1/\alpha - 1 + 2\alpha - 1/(\alpha - 1)^2 \log(\alpha))^2 \alpha b_0 \log^2 S
\]

Figure 6 shows (5), normalized with respect to \( \log^2 S \), as a function of \( \alpha \) for different values of \( b_0 \). One notices that for increasing bit sizes the branching ratio of the hierarchy should decrease.

Because of the simplicity of their storage cells, dynamic memories have \( b_0 \) between 1 and 2. Static memories require a cross-coupled structure and hence a larger \( b_0 \)--typically 3 to 4. For optimal designs, static memories should therefore have a smaller \( \alpha \) than dynamic ones. For dynamic MOS memories the optimal choice for \( \alpha \) lies between 8 and 16, for static MOS memories between 4 and 8. “Smart memories”, structures in which part of the processing task is distributed over the memory cells, have quite large level 0 modules containing an entire processor. They should therefore have small branching ratios and hence relatively deep hierarchies. Current commercial memory chips are designed with \( \alpha \approx 100 \) at the lowest level. This value approximately minimizes the product of the access time and the exponential of the area. Designs of this sort reflect the near exponential dependence of yield on chip area in the early, low-yield phase of a device’s production history. However, near its production peak, the area-time product is closer to a realistic cost function. This shift in production economics suggests that redesigns of high-volume devices should be done using smaller values of \( \alpha \) than initial designs.

**Energy per Access**

In real systems, the cost of power, cooling, and electrical bypassing often exceeds the cost of the chips themselves. Hence any discussion of the cost of computation must include the energy cost of individual steps of the computation process. In a RAM, each access costs an energy proportional to the length of the wires which must be charged or discharged during a given cycle.

Consider a RAM such as that shown in Figure 4. At the highest level (level N) such a device has \( S = \alpha^{2N} \) bits. In each cycle \( \log S \) address wires of length \( L_N \) will in general change state. In addition one horizontal select line, \( \alpha \) vertical data lines, and one multiprocessor output line (all of
length \( L_N \) will change state. Thus at level \( N \), the energy expended per access will be

\[
E_N \sim L_N \left[ \log S + \alpha + 2 \right]
\]

At level \( N-1 \), \( 2\log \alpha \) fewer address wires will be needed. Since only one select line will be active, only \( \alpha \) of the \( \alpha^2 \) submodules will be active. Each submodule contains wires approximately \( 1/\alpha \) as large as those at level \( N \).

Thus the total energy per access is

\[
E_T \sim L_N \left[ \log S(1 + 1 - 2\log \alpha/\log S + 1 - 4\log \alpha/\log S + \ldots) + \alpha + 2 \right]
\]

This expression evaluates to

\[
E_T \sim L_N \log S/\log \alpha \left[ \log S/4 + (\alpha + 2)/2 \right]
\]

Using the \( b_0 \) values from (4), the energy per access of any given size RAM may be evaluated. The results of such an evaluation for a 65K bit RAM are shown in Figure 7.

These curves suggest that considerably less power would be required if memory chips, even of current size, were built with smaller submodules and smaller \( \alpha \).

General Method of Analysis

We have presented a general method for analyzing the cost and performance of recursively defined VLSI structures. Parameters of any such structure may be optimized with respect to some combination of access time, area, and energy.

The results of this study indicate that as more processing is available in each module at level zero, \( b_0 \) will be larger and the optimal value of \( \alpha \) will decrease. A system with \( \alpha = 4 \) would seem to be appropriate for structures in which substantial processing is commingled with storage.

Very general arguments were used to generate the basic recursive structure. For that reason it appears that a very large fraction of VLSI computing structures will be designed in this way. The way in which the area, time, and energy measures were established should make it clear how to apply these techniques to other recursively defined computing structures.
Fig. 7   Energy per Access as a function of Alpha
5. Challenges for the Future

We have seen that it is possible to construct general-purpose computing engines that exploit tremendous concurrency if computations are properly matched to the machine. The vast quantity of concurrency available in such machines can be an enormous help with the computing tasks we face. However, to date we have no formal way of making the possible concurrency in any given calculation apparent or finding if we have come close to the possible concurrency inherent in the computation.

The future of concurrent processing is bounded in part by our ability to escape the strong hold that the conventional sequential machine exerts on our thinking. We must approach problems with concurrency in mind, recognizing that communication is expensive and that computation is not. Progress in these endeavors will surely increase when some VLSI computers of the sort we have illustrated in this chapter begin to appear. When the effort of casting the problem as a structure of concurrent processes is rewarded by a tangible increase in performance, the incentive to design concurrent algorithms will surely increase.

The tools that we use to design and implement concurrent processes are primitive. We are badly in need of notation or language that expresses the power and constraints of highly concurrent machines. Whether such machines are general- or special-purpose, a natural way is needed to map problems onto them. Only in this way will it be possible for applications to find their way into execution in this new computing environment rapidly. In addition we need a method of formally proving the correctness of algorithms mapped onto such machines; it is not possible for human programmers to keep track of the exact relationship of the enormous number of tasks executing on such a machine. An ideal notation would allow expression of only those operations which are free of obvious fatal errors such as deadlock. Only one such notation is known to the authors at this writing, that of the *Association* by Martin Rem [15].

Perhaps the greatest challenge that VLSI presents to computer science is that of developing a theory of computation that accommodates a more general model of the costs involved in computing. Can we find a way to express computations that is independent of the relative costs of processing and communication, and then use the cost properties of a piece of hardware to derive the proper program or programs? The current VLSI revolution has revealed the weaknesses of a theory too solidly attached to the cost properties of the sequential machine.
References

1. I.E. Sutherland and C.A. Mead, "Microelectronics and Computer Science," *Scientific American*, vol. 237, no. 9, September 1977, pp. 210-228. This article is a readable and inspiring call to abandon conventional computing structures.


Chapter 9: Physics of Computational Systems

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Sections:

-Thermodynamic View of Computation
- Energetics of Bistable Devices
- Thermal Limit
- Quantum Limits
- Granularity of Charge
- Voltage Limit
- An Example
- Energy Management
- Discreteness in Quantum Mechanical Systems
- Conclusion

Computation is in the end a physical process. Data elements must be represented by some physical quantity in a physical structure, for example, as the charge on a capacitor or the magnetic flux in a superconducting ring. These physical quantities must be stored, sensed, and logically combined by the elementary devices of any technology out of which we build computing machinery. At any given point in the evolution of a technology, the smallest logic devices have a definite physical extent, require certain minimum time to perform their function, and dissipate a switching energy when switching from one logical state to another. From the system viewpoint, these quantities are the units of cost of computation. They set the scale factor on the size, speed, and power requirements of a computing system. Some of the relationships between these elementary quantities are discussed in this chapter, and an example is given of application to technology comparison.

Thermodynamic View of Computation

{ in preparation }

Energetics of Bistable Devices

Any physical structure we use to represent information must be reliable. We must be able to stably store all bits of information in our computing machine over the period of any computation. Binary information implies elementary memory elements of a bistable nature: one state denoting a logical zero, the other a logical one. A mechanical system which behaves in this way is the inverted pendulum shown in figure 1a. The force of gravity holds the pendulum stably in either the rightmost or the leftmost position. Switching from one state to the other can be accomplished by pushing the weight up to its maximum position and letting it fall onto the opposite stop.

Physicists view bistable systems of this sort in terms of a diagram such as that shown in figure 1b. What is plotted here is the potential energy of the physical system as a function of its spatial or
electrical coordinate. If the pendulum is left in one of its stable states, given by the minima in the potential diagram, it will stay there indefinitely until enough external energy is provided to surmount the potential maximum and allow the system to re-equilibrate in the other potential minimum. Note that the energy provided by the external switching source is lost in the impact when the pendulum falls to its stop (and perhaps bounces a bit until the energy is dissipated). Others have considered particles in potential wells of this shape to derive minimum switching energies for computation.\textsuperscript{5}

The slope of the energy curve, i.e. the derivative of the energy with respect to the angle of the pendulum, has the units of a torque. This torque is being supplied by gravity, and pushes the pendulum towards one of its stable positions. Note that gravity acts as the "power supply" for this mechanical logic device.

The energy required to switch from one state to the other can be supplied deliberately, or by some random occurrence. Suppose our pendulum were mounted on a railroad car. While the train is stationary, we expect the device to remain in its initial state. However when the train passes over a very rough stretch of track, the pendulum may bounce into the other state. The potential maximum must be high enough to prevent such random events.

An electronic circuit with the same logical behavior as the pendulum is the ordinary flip-flop shown in figure 2. The detailed behavior of the flip-flop is, however, somewhat different from that of the inverted pendulum. Let us attempt to change the state of the device by supplying a current into that side of the flip-flop which is at the lower potential. If the current we supply is large enough, we will raise the potential on that side of the flip-flop, turn on the transistor on the opposite side, and change the state of the flip-flop. We can, however, supply a lower current (and therefore power) for an indefinite period of time without changing the state of the device.

In the pendulum we could support the weight part way up the potential curve for a long time and not change the pendulum’s state. However, while supporting the weight in a fixed position, we would not be supplying power to the pendulum. We supply power to the pendulum only when we are increasing the elevation of its mass in the gravitational field. It is thus clear that while in general the behavior of the flip-flop and the inverted pendulum are similar, the detailed energetics are quite different. In particular, we can supply a large quantity of energy to the flip-flop without changing its state, provided we supply the energy slowly enough. This is not true of a system like the inverted pendulum.
Fig. 1a Inverted Pendulum

Fig. 1b Potential energy of Pendulum
Fig. 2 The Energetics of a Flip-Flop
Principle of Least Power

The basic physical law governing the behavior of any electrical circuit involving resistors is called the Principle of Least Power\(^1\). Any electrical network composed of resistors comes to equilibrium by adjusting the potentials in such a way that the power dissipated in the network is a minimum. This principle holds true even if the resistors which form the network are not ideal linear resistors. Any network composed of dissipative electrical elements, such as the MOS or Bipolar transistors, will behave in this way.

Energetics of the Flip-Flop

The power dissipated by our ordinary MOS flip-flop when we forcibly hold one node (\(V_1\) for example) at an arbitrary voltage, is plotted in figure 2. Notice that the curve has the same general shape as the energy curve for our inverted pendulum. The two minima correspond to the two stable states of the flip-flop. The maximum corresponds to the point at which no external power need be supplied to hold the flip-flop in its intermediate state, i.e. it is the metastably balanced condition for the flip-flop.

The derivative of the total power with respect to the voltage \(V_1\) has the dimensions of a current, and is equal to twice the current we must force into the node to hold it at a particular voltage.

Although the principles we derive for circuits of this sort are quite general in nature, it is instructive to work a simple, idealized example. Let us represent the pullup transistors of our flip-flop as ordinary resistors with resistance \(R\), and the pulldown transistors as current sources whose magnitude is some mutual conductance \(G_m\) multiplied by the voltage above threshold of the transistor gate.

This idealized equivalent circuit is shown in figure 3a. The transfer characteristic of each individual inverter in the flip-flop is shown in figure 3b. The output voltage of the inverter is constant at voltage zero until the input voltage exceeds the threshold voltage \(V_{th}\) of its pull down transistor. \(V_{out}\) then varies linearly with a slope \(-G_mR\). This slope is the gain of the inverter. After the output voltage reaches 0, the inverter saturates and its output remains 0 for further increases in the input voltage.

The power dissipated by this circuit for any given value of the voltage \(V_1\) can be computed analytically and is plotted in figure 4 for various values of the transistor transconductance \(G_m\).
Notice that for values of $G_mR > 1$ the power curve shows a distinct maximum in the center separating the two minima corresponding to the stable states of the flip-flop. However, when the gain is one or less, the power curve shows only one minimum near the threshold voltage. This minimum corresponds to a single stable state.

*A cross-coupled circuit must have a loop gain greater than unity, in order to develop two independent stable states.* [anon]

This analysis based on the Principle of Least Power agrees with this standard electrical engineering model. It provides us with a very general and fundamental viewpoint from which to analyze the energetics of computer circuits.

Let us perform a conceptual experiment on the flip-flop of figure 2. In its initial state, $V_1$ is a logical one and $V_2$ is a logical zero. We delicately remove the connection from $V_2$ to the gate of the left transistor. At some instant of time we place a charge equivalent to a logical one on the (now floating) gate. At the identical instant we force $V_2$ to a logical one.

As time progresses we observe the power we must supply to keep $V_2$ at a logical one. For a while the device absorbs a large amount of power. However, after the signal from the gate has propagated through the two inverters, no more power will be absorbed, and we may reconnect the gate to $V_2$. Try as we might, we can find no path from one state to the other which can be traversed without supplying an amount of power at least as high as the maximum in the power curve. The total energy required is at least the product of this power and the inverter pair delay.

The essential difference between the "static" and "dynamic" storage devices discussed in chapter 7 is thus clear. Both forms use the same physical element to store the energy which represents information. However the "dynamic" form requires only that the requisite amount of energy be supplied to change its state. It does not matter how slowly that energy is supplied. The "static" form requires in addition that the energy be supplied within the inverter pair delay of the technology.
Fig. 3a  Equivalent Circuit

Fig. 3b  Inverter Transfer Characteristic
Fig. 4  Power curves for Flip-Flop of fig. 3
Thermal limit

We have just illustrated how to compute the energy required for switching the flip-flop. An external influence must supply an additional power, $P$, equal to the difference between the power at the minimum and maximum of the power curve in order to switch the flip-flop from one stable state to the other. That power must be supplied long enough for the information to propagate through both inverters, and back to the node where we applied the signal. This time is just the inverter pair delay $\tau$ for the technology out of which the flip-flop is built. The externally input energy required to flip the flip-flop thus becomes:

$$E_{sw} = \tau P$$

The switching energy must be sufficient to prevent random occurrences from changing the state of the device. Electrical noise is always present in any real system. It is generated by heavy electrical equipment and propagated along power mains. Radio and television transmitters of all varieties create electromagnetic radiation which can induce voltages in a circuit. Modern electronic devices allow single electronic occurrences to control relatively large currents. Atomic imperfections randomly capture and release electrons, thus creating an unsteady environment. Techniques exist for minimizing the effect of each such hazard. However, one fundamental source of irreducible randomness remains.

Any device operating at a finite temperature is subject to the random thermal motions of the elements of which it is composed. The energy of any element, large or small, is not fixed, but fluctuates over a range of energies due to interactions with its environment. Each time we measure the energy $E$ of an element, it will have some value which differs by some $\Delta E$ from its equilibrium value. The probability that any given independent measurement will yield a given $\Delta E$ is given by the Boltzmann equation: $\text{Probability} = e^{-\Delta E/kT}$.

What constitutes "independent" measurements depends on the response time $\tau$ of the system. Two measurements should be made at least $\tau$ apart to be considered independent. Similarly we have seen that in order to switch a bistable system from one state to the other, a certain amount of power $P$ has to be supplied for the duration of the response, or switching, time $\tau$ of the system. Therefore, systems with a faster response time are more likely to be switched by thermal fluctuations, since occasions where the critical power level $P$ is exceeded for the necessary amount of time $\tau$ occur more often. This is equivalent to the view that systems with a wider bandwidth capture more energy out of the spectrum of the thermal energy. Thus the probability per unit
time that an element will achieve some large variation \( \Delta E \) is:

\[
\text{Probability per unit time} = (1/\tau)e^{-\Delta E/kT}
\]

The probability per unit time that random thermal noise will change the state of a bistable device is thus:

\[
\text{Probability per unit time} = (1/\tau)e^{-E_{sw}/kT}
\]

Typical computations may involve many millions of individual memory elements over a period of many hours. Hence we must insist that the probability of spontaneous switching be less than one part in \( 10^{12} \) or so, which requires a switching energy energy of the order of \( 30kT \).

Viewing these energetic considerations from a system level, we have established an absolute minimum for the energy required for doing any given computation.

*The energy required for a computation has a lower bound given by the minimum switching energy multiplied by the number of elementary switching events which must occur during the computation.*

This estimate of minimum energy completely ignores the energy cost of communicating data from one location to another. In many systems, the total communication energy is much larger than the total switching energy.

In realizable electronic systems, the switching cost for elementary storage elements is much larger than the limit given above. In typical 1978 MOS technology, switching an elementary flip-flop requires \( 10^{-12} \) Joule, or approximately \( 10^8kT \) at room temperature. Even a 1/4 micron MOS transistor will require approximately \( 10^4kT \); more than 100 times the energy necessary for reliable computation as given above.

Note that this view of computation makes it perfectly clear that there is no possibility of 100% reliable computing systems. There is always a finite chance that some storage element will switch spontaneously due to thermal noise. However, in today's systems, and even foreseeable VLSI systems, the probability of such a random switching event due to thermal noise is much less than that of a failure due to electrical noise, cosmic rays, or mundane device failure mechanisms. In systems with poorly designed timing constraints, synchronizer failures occur many orders of magnitude more frequently than thermal failures. This observation is the origin of the *Seitz Criterion* given in Chapter 7.
Quantum Limits

The thermal limit given above represents only one way in which an immutable law of nature places bounds on what can be physically realized. Other such limitations come from other physical laws. The lower bound on the size of an FET which will operate properly is determined, not by thermal considerations but by the uncertainty principle and the discreteness of electrical charge. From the uncertainty principle, an electron of mass \( m \) will, because of its wave nature, have an uncertainty \( \Delta x \) in its position \( x \) related to the uncertainty \( \Delta p \) in its momentum \( p \) by:

\[
\Delta p \Delta x \approx \hbar
\]

The energy is related to the momentum by:

\[
E = \frac{p^2}{2m}
\]

Hence an energy barrier of thickness \( \delta \) and height \( E_b \) can contain an electron only if:

\[
\delta \gg \Delta x \approx \frac{\hbar}{(2mE_b)^{1/2}}
\]

For a barrier of height 1eV, \( \Delta x \) is about 0.001 micron. Gate oxides and junction depletion layers must be many times this thickness. In 1978, gate oxide is already less than 0.1 micron thick. We are thus within sight of a fundamental size limitation due to quantum phenomena.

Granularity of Charge

An even more severe limit results from the discreteness of impurity ion charges in the depletion layer under the FET channel. Let us attempt to reduce all voltages \( V \) and distances \( d \) by the same factor. A charge layer of \( q \) charges per unit area produces an electric field. This field across a depletion layer of thickness \( d \) results in a voltage \( V \):

\[
V \propto qd
\]

The charge is due to impurity ions of density \( N \) per unit volume. Hence:

\[
q \propto Nd
\]

The voltage is therefore proportional to the square of the depletion layer thickness:

\[
V \propto Nd^2
\]
In order to scale both $V$ and $d$ by the same factor, $N$ must therefore be proportional to $1/d$.

The total number of charges in the channel is the number per unit volume times the volume of the region under the gate. By our scaling convention, all volumes must be proportional to $d^3$:

$$N_{\text{tot}} \propto N d^3 \propto d^2$$

For randomly distributed impurities, the expected statistical variation of the total number $N_{\text{tot}}$ is:

$$\Delta N_{\text{tot}} = (N_{\text{tot}})^{1/2} \propto d$$

Or a fractional variation of:

$$\Delta N_{\text{tot}}/N_{\text{tot}} \propto 1/d$$

This statistical variation of the number of impurity ions under the channels of different transistors results in a similar distribution of threshold voltages:

$$\Delta V_{\text{th}}/V_{\text{th}} = \Delta N_{\text{tot}}/N_{\text{tot}} \propto 1/d$$

The variation in threshold voltages thus becomes larger as devices become smaller. A detailed treatment of this effect is given in [Ref. R3 of Ch.1], which concludes that a device of 4 μm channel length described in [Ref. 3 of Ch.1] would have an expected variation in its threshold voltage of $\approx 0.08$ Volts. There are techniques which can greatly reduce this statistical variation. A thin film of undoped silicon just under the gate oxide will largely isolate the threshold voltage from the granularity of charge in the substrate. Such a structure complicates an already difficult task of sub-micron fabrication. It therefore appears that, aside from the fabrication process, the first barrier we face in the sub-micron FET world is a difficulty in scaling voltages to low enough values. We consider the fundamental limit on supply voltage in the next section.

**Voltage Limit**

We have seen that a storage device must exhibit a maximum in its power curve in order to retain information. There are two independent ways in which this maximum may become too small. The first is that the elementary logic gates may become too small to store enough energy. We see that this limit does not constrain ordinary FET logic since FET gate lengths must be greater than 1/4 micron for other reasons. Another way is that the operating voltage may become too small to assure that the gain of an elementary circuit exceeds unity. Semiconductor technology is evolving
Fig. 5 Conceptual Model of a Complementary Inverter at very low Voltages
under a scaling law in which operating voltage must be decreased along with device dimensions. Hence it is important to establish a lower limit on the operating voltage of FET circuits.

As mentioned in Chapter 1, when a MOS device is operated near its threshold, the channel resistance $R_{ch}$ is exponentially dependent upon the gate voltage $V_g$:

$$R_{ch} \propto e^{-qV_g/nkT}$$

The factor $n$ is due to the substrate effect, and is approximately 1.2 for most processes.

A model of a complementary device (such as a CMOS inverter) is shown in figure 5. The resistances $R_1$ of the lower n-channel device and $R_2$ of the upper p-channel device are exponentially dependent on the input voltage $V_{in}$ as follows:

$$R_1 = R e^{-qV_{in}/nkT}, \quad R_2 = R e^{qV_{in}/nkT}$$

Therefore the output voltage is:

$$V_{out} = \left[ V/2(1-R_1/R_2) \right] [R_1-R_2] V_{in}$$

We are interested in the gain near the switching threshold, which because of the supply voltage convention is at $V_{in} = 0$. We may expand the exponentials as a power series and ignore all but the first order terms in $V_{in}$:

$$V_{out} \approx -(qV/nkT)V_{in}$$

The gain of the circuit is thus equal to $qV/nkT$. Hence realistic supply voltages for complementary circuits should be a few $kT/q$. At room temperature ($kT/q \approx 25$ mV). Ratio logic families, such as nMOS, can be analyzed by the same technique. Since they have only one non-linear device rather than two, their gain is approximately half that given above. They will therefore require twice the supply voltage required by complementary devices. Routine CMOS circuits with $\approx 5$ micron geometries operate with a 5V supply. Scaling in a straightforward way, we would expect $\frac{1}{4}$ micron devices to operate with a $\frac{1}{4}$ Volt supply.

While the gain of such a circuit would be adequate if all its transistors had the same threshold voltages, it is possible that the pullup transistor of an inverter could have a particularly high threshold voltage, while its companion had a particularly low threshold. If the difference in threshold voltages exceeded the supply voltage $V_{pp}$ the device output would always remain in one state. The probability of such an occurrence, computed from the variation in threshold
mentioned in the last section is:

\[ P = e^{-2V_c/\Delta V_{th}} \]

We might for example require that, even in a VLSI system containing \(10^7\) inverters, the probability of all the system's transistors being within threshold limits be greater than 0.9. Such a criterion would require a supply voltage of \(\approx 0.7\)V. Unless special attention is directed toward reducing threshold variations, systems with \(1/4\) micron device geometries will be forced to operate with higher supply voltages than the straightforward scaling would indicate. However the inherent nonlinearity of the FET near threshold lowers the effect of threshold variations, and system operation at a supply voltage in the 100 - 200 mV range appears feasible.

An Example

In this section we will apply physical considerations to the comparison of two very different technologies for constructing computational systems. The technologies selected for this example are based on (i) semiconductor FET devices and (ii) Josephson junction devices. The material presented in this example demonstrates the importance of considering not only device physics and device design, but also system physics and system architecture, when making such comparisons.

Several types of limits on the performance of semiconductor FET logic families have been noted in the foregoing discussion: those dealing with the temperature of operation, those arising from quantum phenomena, those associated with the granularity of charge in the semiconductor substrate, and voltage constraints arising from gain considerations. Of these, the limit due to quantum phenomena appears the least restrictive.

It would thus appear that a physical process not involving a doped semiconductor and operating at very low temperature would merit serious study. Superconducting logic families have, for this reason, attracted much attention. Information is stored as a magnetic flux trapped in a superconducting ring, and is switched by means of a Josephson (or similar) junction. Devices have been demonstrated which exhibit very fast switching times and low operating power. It is important to understand the relative merits of such a radically different technology from the point of view of overall system design. We should therefore find some way to compare it directly with semiconductor technology, and to extend the comparisons to scaling into sub-micron dimensions.

In real systems, the cost of energy, and energy conversion and distribution, often exceeds the cost
of the chips themselves. Hence any discussion of the cost of computation must include the energy cost of individual steps of the computation process. The fundamental figure of merit of a logic device is its *switching energy* discussed previously. This quantity is a measure of the power-delay product of the technology. Propagation delay can be traded off against power dissipation over a wide range in any given technology, but their product cannot be reduced below the switching energy. In a charge controlled semiconductor device such as the MOSFET, the irreducible switching energy is \( E_{sw} = C_g V^2 / 2 \), for gate capacitance \( C_g \) and supply voltage \( V \).

In a superconducting device \( E_{sw} = LI^2 / 2 \), where \( L \) is the inductance of the superconducting loop plus the associated junction, and \( I \) is the supply current. In both technologies, parasitics will increase \( E_{sw} \) to several times the values computed for minimum devices. However, for purposes of comparison, we will consider only the minimum devices themselves.

Since all energies in both types of logic are multiples of \( kT \), it might appear that operating a computer at very low temperatures would reduce the total power required. That this is not the case is easily demonstrated. Suppose that to perform a computation a machine dissipates energy \( E_L = nkT_L \) as heat at some low temperature \( T_L \). To maintain the low temperature, this heat energy must be transported to and released at room temperature, \( T_H \), by some refrigerator. The total energy to run the system is equal to \( E_L \) plus the work required to run the refrigerator. Thermodynamics shows us that a refrigerator operating on the Carnot cycle requires the least amount of work input per unit of heat transported from the low temperature environment to the high temperature environment.\(^3\) On input of work \( W \), a Carnot refrigerator can transport, from the \( T_L \) to \( T_H \) environments, a quantity of heat energy \( Q \) given by:

\[
Q/W = T_L / (T_H - T_L)
\]

Thus the work \( W \) required to transport \( E_L \) from \( T_L \) to \( T_H \) is in general:

\[
W \geq E_L(T_H - T_L) / T_L
\]

The total energy, \( E_{tot} \), required for the computation is therefore:

\[
E_{tot} \geq nkT_L + nkT_L[(T_H - T_L)/T_L] = nkT_H
\]

As \( T_L \) is lowered, the switching energy is lowered, but the work input to the refrigerator must be increased by at least an equal amount. The total energy cost, including that necessary to run the refrigerator, is thus independent of the temperature of the computer's switches. This energy cost
is, at minimum, identically equal to nkT at the temperature of the ultimate heat sink. In some space applications a heat sink at very low temperatures is available. However, for terrestrial computers, refrigerating electronic devices in order to reduce the energy of computation is logically equivalent to constructing a perpetual motion machine. For this reason, we will use kT at the heat sink temperature in system energy calculations, independent of the actual temperature at which the switching devices operate.

Now we turn to the details of the technology comparison. The switching energy of MOSFET logic is: \( E_{sw} = C_g V^2/2 \). The most straightforward MOSFET scaling results from reducing all dimensions by the same scaling factor. If this type of scaling is applied to the MOS family, the gate capacitance decreases linearly with the scaling factor. In order to keep the electric fields constant, the supply voltage is scaled by the same scaling factor. The switching energy is thus reduced by the third power of the scaling factor, as illustrated in the top curve in Figure 6. The lower size limit shown is a conservative estimate set by device physics factors previously discussed.

Were it possible to build FET devices which operated with one electronic charge on their gate, their performance would not benefit from scaling to smaller dimensions. In such a device, the switching energy can be expressed in terms of \( q_0 \), the charge of the electron:

\[
E_{sw} = CV^2/2 = q_0^2/2C
\]

Since C decreases as the device dimensions are scaled down, the switching energy actually increases. This relationship illustrates a general principle: A logic device working at its quantum limit requires a higher switching energy as the dimensions of the device are made smaller.

Even at present dimensions, superconducting logic operates at or near its quantum limit. The flux in a superconducting ring must be an integral multiple of the flux quantum \( \Phi_0 \approx 2 \times 10^{-15} \) Webers. The switching energy for a device operating with one flux quantum can be written as:

\[
E_{sw} = LI^2/2 = \Phi_0^2/2L
\]

Note that the inductance \( L = \Phi_0/1 \) is directly proportional to the size of the loop. The above dependence for superconducting logic is illustrated in the bottom curve in Figure 6. The lower size limit shown is set by the penetration depth \( \lambda \) of the superconductor. Magnetic field strength decreases with distance, \( x \), into the superconductor as \( e^{-x/\lambda} \). If the thickness of the superconducting ring is less than a few \( \lambda \), the ring cannot localize the flux within it. A typical
Fig. 6 Comparison of FET and Superconducting Logic

- **FET Logic**: $E_{sw} = \frac{CV^2}{2}$
- **Superconducting Logic**: $E_{sw} = \frac{\phi_0^2}{2L}$
value of λ is 0.1 micron.

Comparing the upper and lower curves, it is clear that, when an accounting is made of the total energy, and when the effects of scaling to sub-micron dimensions are taken into account, room temperature FET logic is a remarkable technology. At achievable sub-micron dimensions, it can actually outperform its superconducting counterpart. Lower switching energies in the superconductor technology can be achieved only by sacrificing density. This trade-off may be desirable under some circumstances. It seems more likely, however, that maximum computation per unit cost will be achieved by jointly minimizing switching energy and maximizing circuit density.

The absolute speed attainable with the superconducting logic is, however, considerably better than that of its FET counterpart. For a critically damped Josephson junction, the time response \( \tau \) is

\[
\tau \approx (L/C)^{1/2},
\]

where \( L \) is the loop inductance used above, and \( C \) is the junction capacitance. Since the normal resistance of the Josephson junction varies exponentially with dielectric thickness, the thickness can be assumed approximately constant as the devices are scaled. Hence the delay time \( \tau \) will scale down as the 3/2 power of the scaling factor. For the FET, the oxide thickness must be scaled, and the delay time varies linearly with the scaling factor.

At 1 micron feature size, for example, the switching time of a superconducting device is \( \approx 2 \times 10^{-13} \) sec, while for a FET with the same feature size the transit time is \( \approx 10^{-11} \) sec.

One basic problem with low temperature logic is that the lower switching energy levels result in poor noise immunity. They therefore require better shielding to reduce the effect of external electromagnetic occurrences to a level well below the switching energy.

Another problem is that the low switching energy creates a mismatch to the outside world for which a penalty in additional power consumption has to be paid, since the drivers to the outside world consume a large amount of power, and introduce extra delays. As long as information is not required to exit the low temperature environment, chip to chip communication can be done at high bandwidth. Note that in this respect, superconducting logic is superior since it is somewhat better matched to the impedance of transmission lines than is FET logic. In any event, exponentially staged drivers are required when driving from the low energy environment to the
outside world, as discussed in chapter 1. These drivers introduce a minimum delay $\tau_{dr}$:

$$\tau_{dr} \geq re \ln(Y).$$

where $Y$ is the ratio of energy required at the destination to that of the elementary logic device. If the switching energy of a logic element is a factor of 100 smaller due to operation at low temperature, a factor of at least 10 in driver delay is introduced. Furthermore, the dissipation of the last stage of the driver is determined by the energy level necessary in the outside world, not in the low temperature environment. The cost of this driving energy is at least 100 times higher than that for a room temperature driver of the same capability, due to the constraints imposed by the laws of thermodynamics.

It is important to recognize that the trade-off between power and delay time extends to much shorter times for Josephson devices than than it does for FET's. The speed advantage of the Josephson devices, in the scaled environment of the future, will be about a factor of fifty. Although their switching energy will be about the same as that of FETs, we would have the option of inputting fifty times more power into a system composed of Josephson devices, and then being able to switch them fifty times faster than the fastest FETs.

Architects comparing alternative technologies for building computing systems take into account many costs other than just total switching energy. The weights assigned to the various factors usually depend upon their proximity to absolute constraints imposed by physical law or by system performance and cost considerations. In certain situations, we may be perfectly willing to pay the price for large increments in energy, energy conversion equipment, mass, volume, and structural and operational complexity, in order to achieve an increment of system performance.

Suppose, for example, we now had to specify a very high performance general purpose computer for the late 80's or early 90's. Since switching speed translates directly into time performance in the classical stored program computer, we might see no other alternative for high performance than a machine based on superconducting devices. Such a decision recognizes that no present alternatives exist for trading off processing speed against concurrency in multiple processors for general purpose computation. That such alternatives must ultimately exist is of course evident by observation of the information processing capability of living organisms.

Superconducting devices meet the requirement for high speed in the classical computer, and a number of machines based on that technology will likely be built before viable high concurrency
alternatives appear. However, in the longer term, in applications where mass, volume, structural complexity, and cost are real constraints, semiconductor devices operated at heat sink temperature will generally have the advantage. Thus, the switching technology likely to dominate the terrestrial environment, used for personal computing and personal communications on a vast scale in an enormous number of different applications, is semiconductor technology. Recall that semiconductor technology itself may benefit in a variety of ways from low temperature operation [Ref.7 of Ch.1], as for example in the reduction of subthreshold current in submicron MOSFETs.

Energy Management
{ in preparation }

Discreteness in Quantum Mechanical Systems
{ in preparation }

Conclusion

We opened this book with a discussion of the physical properties of elementary switching devices. We have now closed with a discussion of fundamental physical principles which profoundly influence the higher-level properties of computing systems.

The communication of information over space and time, the storage of information by change of state at storage sites, and the transport of energy into and heat out of systems depend not only on abstract mathematical principles, but also on physical laws. The generation and synthesis of very large scale systems, whether artificial or natural, proceed under and indeed are directed by the constraints imposed by the laws of physics.

We look forward to a time when quantitative measures can be given for the true cost and complexity of any required computation. At present we are very far from this goal. The examples of this chapter do, however, serve to illustrate that concrete physical arguments can be applied to the properties of information systems. We hope others will provide insights and examples in this important area of investigation, for reporting in future editions of this text.
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