01 Jan 1971

# A State Assignment Procedure For Asynchronous Sequential Circuits 

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## Recommended Citation

G. K. Maki and J. H. Tracey, "A State Assignment Procedure For Asynchronous Sequential Circuits," IEEE Transactions on Computers, vol. C thru 20, no. 6, pp. 666-668, Institute of Electrical and Electronics Engineers, Jan 1971.
The definitive version is available at https://doi.org/10.1109/T-C.1971.223320

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$$
\left[\begin{array}{llll}
0 & 1 & 0 & \cdots
\end{array}\right) 00
$$

representing the cyclic permutation

$$
(012 \cdots(m-1)) \text { on }\{0,1, \cdots, m-1\} .
$$

Recall the definition of the Kronecker product of two matrices [7]. If $A$ is an $m \times m$ matrix, the Kronecker or tensor product $A \otimes B$ is the $m n \times m n$ matrix

$$
\left[\begin{array}{ccc}
a_{11} B \cdots & & \\
& \ddots & a_{1 m} B \\
\vdots & & a_{i j} B \\
& & \ddots \\
a_{m 1} B & \cdots & \\
& & a_{m m} B
\end{array}\right]
$$

The Kronecker product is associative. Also, it is easy to see that the Kronecker product of permutation matrices is a permutation matrix.

Furthermore, the permutation represented by $Q_{n_{1}} \otimes Q_{n_{2}}$ $\otimes \cdots \otimes Q_{n_{k}}$ can be described in "coordinatized" form as follows. It maps

$$
j_{1}+j_{2} n_{1}+\cdots+j_{k} n_{1} \cdots n_{k-1}
$$

into

$$
\left(j_{1}+1\right)+\left(j_{2}+1\right) n_{1}+\cdots+\left(j_{k}+1\right) n_{1} \cdots n_{k-1}
$$

where $j_{p}+1$ means addition modulo $n_{p}$. It is then straightforward to verify that

$$
\begin{equation*}
P^{-1} Q_{n} P=Q_{n_{1}} \otimes Q_{n_{2}} \otimes \cdots \otimes Q_{n_{k}} \tag{2}
\end{equation*}
$$

Let $A(x)$ be an $n \times n$ circulant. Then

$$
A(\boldsymbol{x})=\sum_{j=0}^{n-1} x_{j} Q_{n}^{j}
$$

where $Q_{n}^{0}=I_{n}$, the $n \times n$ identity matrix. Thus, applying (2),

$$
\begin{align*}
P^{-1} A(x) P & =\sum_{j=0}^{n-1} x_{j}\left(Q_{n_{1}} \otimes \cdots \otimes Q_{n_{k}}\right)^{j} \\
& =\sum_{j=0}^{n-1} x_{j}\left(Q_{n_{1}}^{j} \otimes \cdots \otimes Q_{n_{k}}^{j}\right)  \tag{3}\\
& =\sum_{j=0}^{n-1} x_{j} Q_{n_{1}}^{j_{1}} \otimes \cdots \otimes Q_{n_{k}}^{j_{k}} \tag{4}
\end{align*}
$$

Equation (3) follows from the matrix identity $(A \otimes B) \cdot(C \otimes D)$ $=(A \cdot C) \otimes(B \cdot D)$, while (4) follows from the identity $Q_{p}^{p}=I_{p}$ for all $p$. If $C_{i}$ is an $n_{i} \times n_{i}$ circulant for $i=1,2, \cdots, k$, then $C_{1} \otimes \cdots \otimes C_{k}$ is an $\left(n_{1}, \cdots, n_{k}\right)$ supercirculant. Finally, a linear combination of $\left(n_{1}, \cdots, n_{k}\right)$ supercirculants is an $\left(n_{1}, \cdots, n_{k}\right)$ supercirculant. Thus $P^{-1} A(\boldsymbol{x}) P$ is an $\left(n_{1}, \cdots, n_{k}\right)$ supercirculant and the lemma is proved. A more conceptual proof appears in [5]. As a consequence of the supercirculant lemma we obtain the following.

Speedup Lemma: Suppose there is a function $f: N \rightarrow N$, where $N$ is the set of positive integers such that for any
commutative ring $R$, the product of two $n \times n$ circulants can be computed with $f(n)$ multiplications. Then, if $n=n_{1} \cdots n_{k}$, with the $n_{i}$ 's relatively prime in pairs, the product of two $n \times n$ circulants can be computed with $f\left(n_{1}\right) \cdots f\left(n_{k}\right)$ multiplications.

Proof: By the supercirculant lemma it suffices to consider multiplication of two $\left(n_{1}, \cdots, n_{k}\right)$ supercirculants. The proof is by induction on $k$. The assertion is trivially true for $k=1$. Assume that it is true for $k$ and let $S_{1}, S_{2}$ be two $\left(n_{1}, \cdots, n_{k}, n_{k+1}\right)$ supercirculants. Let $R_{k}$ be the set of all ( $n_{1}, \cdots, n_{k}$ ) supercirculants over $R$. It is easy to see that $R_{k}$ is a commutative ring, under matrix addition and multiplication. $S_{1}$ and $S_{2}$ can be considered $n_{k+1} \times n_{k+1}$ circulants over $R_{k}$. Thus $S_{1} S_{2}$ can be computed using $f\left(n_{k+1}\right)$ multiplications in $R_{k}$. Further, by the induction hypothesis each multiplication in $R_{k}$ requires $f\left(n_{1}\right) \cdots f\left(n_{k}\right)$ scalar multiplications. Thus the total number of scalar multiplications required is $f\left(n_{1}\right) \cdots f\left(n_{k}\right) f\left(n_{k+1}\right)$. This proves the lemma.

By Lemma 1, we can take $f(n)=n(n+1) / 2$; this we get the following.

Proposition: Let $n=n_{1} \cdots n_{k}$ with $\left(n_{i}, n_{j}\right)=1$ for $i \neq j$. Then the product of two $n \times n$ circulants and thus the convolution of two $n$-vectors can be computed using $n\left(n_{1}+1\right)$ $\cdots\left(n_{k}+1\right) / 2^{k}$ scalar multiplications.

Remark: It is easy to see that the factorization minimizing the number of multiplications by our method is the complete factorization of $n$ into prime-power factors.

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## A State Assignment Procedure for Asynchronous Sequential Circuits

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Abstract-This paper presents a new procedure for constructing nonuniversal shared-row internal state assignments for asynchronous sequential circuits. The method consists basically of establishing an initial code with the minimum number of variables required to dis-

Manuscript received September 12, 1969; revised October 1, 1970, December 2, 1970, and January 13, 1971. This work was supported in part by the National Science Foundation under Grants GK-2017 and GK20190.
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tinguish the states. Then variables are systematically added until a satisfactory assignment is obtained. State variables added in the expansion of an assignment are merely the ExClusive or of state variables in the original assignment. This simple construction procedure terminates with a maximum of $\boldsymbol{m}+[\boldsymbol{m} / 2]$ state variables for a $\mathbf{2}^{\boldsymbol{m}}$-row flow table.

Index Terms-Asynchronous sequential circuits, internal state assignment, nonuniversal assignment techniques, shared-row assignments.

## I. Introduction

The operation of an asynchronous sequential circuit can be described by means of a flow table (see Fig. 1). An important step in the synthesis procedure for these circuits consists of coding the internal states such that all state transitions can be accomplished without the introduction of critical races [1]. The internal state variables, $y_{1}, y_{2}, \cdots, y_{n}$, and corresponding code are called the internal state assignment. The type of assignment considered here is one in which a single $y$ state is assigned to each row of the flow table and remaining $y$ states are used as necessary to complete transitions between nonadjacent rows. The result of such an assignment is that, after an input change, the circuit may sequence through several internal states before becoming stable. Unger refers to these assignments as sharedrow state assignments [2]. This should be contrasted with single transition time assignments where all internal state transitions go directly from unstable states to terminal stable states; no sequencing is permitted. Generally, fewer state variables are required for shared-row assignments.

A state assignment is said to be universal if it is valid for any flow table with a given number of rows; otherwise, it is nonuniversal. A nonuniversal assignment for a particular flow table frequently requires fewer state variables than a universal assignment. Procedures are well established for generating nonuniversal single transition time assignments for asynchronous sequential circuits [3], [4]. This note treats the more difficult problem of generating nonuniversal shared-row assignments [5]-[7].

The assignment procedure presented here does not guarantee a minimum-variable assignment but terminates with $m+[m / 2]$ state variables for a $2^{m}$-row flow table, where

$$
[m / 2] \triangleq\left\{\begin{array}{l}
m / 2, \quad m \text { even } \\
(m-1) / 2, \quad m \text { odd }
\end{array}\right.
$$

For flow tables where the number of rows are equal to or less than $2^{m}$ (and greater than $2^{m-1}$ ), the method consists of first establishing an initial code of $m$ state variables. Additional state variables are obtained, as necessary, by ExclusiveORING variables from the initial code.

It has not been proven that $m+[m / 2]$ state variables are always sufficient for the coding of a $2^{m}$-row flow table by the procedures of this note. However, a search for counterexamples has been unsuccessful.

## II. Development of the Assignment Procedure

A $k$-set of a flow table column consists of $k-1$ unstable states leading to the same stable state, together with that stable state. A column partition is a collection of the $k$-sets that appear in a column of a flow table. The column partitions for the flow table in Fig. 1 are


Fig. 1. Flow table.

$$
\begin{aligned}
\alpha_{1} & =\{1,2,3 ; 4,5\} \\
\alpha_{2} & =\{1,3 ; 2,4,5\} \\
\alpha_{3} & =\{2,3,4 ; 1,5\} .
\end{aligned}
$$

If a transition is to be effected between a pair of states $S_{a}$ and $S_{b}$, then the pair of states $S_{a}, S_{b}$ is called a transition pair. A transition path is the set of $n$-tuples that the circuit assumes in undergoing a transition between states of a transition pair.

The following definitions are useful in the discussion of the assignment procedure.

Definition: A set of state variables $\left\{y_{1}, y_{2}, \cdots, y_{k}\right\}$ is an independent set of state variables if no $y_{i}$ in the set is equal to the $\bmod 2$ sum $\oplus$ of a subset of the other state variables of the set.

Definition: Consider a set of two independent and distinct state variables $\left\{y_{\xi}, y_{\delta}\right\}$. If $y_{\xi} \oplus y_{\delta}=y_{\eta}$, then the complete set $\left\{y_{\xi}, y_{\delta}, y_{\eta}\right\}$ is called a parity set. State variable $y_{\eta}$ is generated from the independent state variables and is therefore called dependent.

## III. Statement and Illustration of the Assignment

A valid set of transition paths for a column partition must not have any $n$-tuples that are common to different $k$-sets. If the column partition consists of only two-sets, then there is a maximum number of transition paths with no elements in common. Therefore, it appears that the most difficult type of flow table to work with is one composed of column partitions with only two-sets. The assignment procedure will be discussed and illustrated for this case first. It will be assumed that the number of flow table rows is between $2^{m-1}$ and $2^{m}$. Next, the procedure will be generalized to the case where $k$-sets are not restricted to two-sets. Details regarding the rationale for each step can be found elsewhere [6].

Step 1: Code the states with $m$ state variables $y_{1}$, $y_{2}, \cdots, y_{m}$ in such a manner that a minimum number of transition pairs are a maximum distance apart. In general, best results are obtained if the states of transition pairs are assigned codes that are close together.

Step 2: Count the number of times each state variable $y_{i}, i=1,2, \cdots, m$, must be excited to effect a distance two or greater transition. This forms the count list.

Step 3: Let $y_{j}$ and $y_{k}$ be two highest value state variables in the count list that are not yet in a parity set. Generate a new state variable $y_{1}$ such that $y_{1}=y_{j} \oplus y_{k}$. Then $y_{1}$, together with $y_{j}$ and $y_{k}$, forms a parity set.

Step 4: Attempt to find the transition paths for each column partition. If a satisfactory set of transition paths exists, then the assignment procedure is complete. If not,
then return to Step 3 and generate another state variable. The process continues until a set of transition paths is found.

The following is an example to illustrate the procedure. Consider the six following column partitions which correspond to a flow table called machine A:

$$
\begin{aligned}
& \alpha_{1}=\{0,9 ; 7,13 ; 3,1 ; 10,4 ; 2,8 ; 11,15 ; 6,14 ; 5,12\} \\
& \alpha_{2}=\{0,6 ; 2,7 ; 3,13 ; 9,1 ; 10,15 ; 4,12 ; 14,8 ; 11,5\} \\
& \alpha_{3}=\{9,4 ; 7,10 ; 0,12 ; 3,2 ; 15,6 ; 5,13 ; 8,14 ; 1,11\} \\
& \alpha_{4}=\{0,14 ; 7,15 ; 9,12 ; 3,5 ; 4,6 ; 11,13 ; 10,8 ; 1,2\} \\
& \alpha_{5}=\{7,10 ; 2,4 ; 15,12 ; 0,6 ; 3,13 ; 1,9 ; 8,14 ; 11,5\} \\
& \alpha_{6}=\{0,13 ; 7,15 ; 3,5 ; 9,14 ; 4,8 ; 1,11 ; 10,12 ; 6,2\} .
\end{aligned}
$$

The codes assigned to the states of the $m$-cube (four-cube) are shown in Fig. 2. It is possible in this case to avoid coding any transition pair with codes a maximum distance apart.

Next, the number of times each state variable is excited, in a distance two or greater transition, are tabulated as follows.

| State Variable | Number of Times Excited |
| :---: | :---: |
| $y_{1}$ | 14 |
| $y_{2}$ | 15 |
| $y_{3}$ | 12 |
| $y_{4}$ | 16 |

On the basis of this count list, a new state variable $y_{5}$ is generated from state variables $y_{2}$ and $y_{4}$ with $y_{5}=y_{2} \oplus y_{4}$. The Karnaugh map representation of this assignment is shown in Fig. 3.

Transition paths can be constructed within each column, and therefore the assignment is satisfactory. The amount of effort involved in finding this state assignment appears to be significantly less than other known techniques.

It is interesting to note that the assignment procedure just described and illustrated always produces the universal four-variable assignment given by Saucier [7].

For the more general case where $k$-sets are not restricted to two-sets, the above algorithm can be easily modified. Each $k$-set is converted to a set of transition pairs. For example, consider a three-set consisting of states 1,2 , and 3 , with state 1 as the stable state. There are three possible sets of transition pairs for a three-set: $(1,2)(1,3) ;(1,2)(2,3)$; $(1,3)(2,3)$; where the parentheses denote the transition pair. It is apparent that more flexibility is associated with $k$-sets where $k$ is greater than two. To account for this, in the initial coding attention is given primarily to keeping the codes of the two-sets as close together as possible. Little effort is expended in minimizing the distance between states of $k$-sets greater than two-sets. Transition pairs that consist of states that are adjacent are most desirable in that the resulting transition path contains just the two states. The means of determining which independent state variables are to be used in generating the additional state variables is identical to the procedure previously discussed for the all-two-set case.

The final code depends on the flow table characteristics and the initial encoding. Therefore, when all transitions

| $y_{1} y_{2}$ |  |  |  |  |
| ---: | :---: | :---: | :---: | :---: |
| $y_{3} y_{4}$ | 00 | 01 | 11 | 10 |
| 00 | 0 | 9 | 13 | 14 |
| 01 | 6 | 11 | 1 | 2 |
| 11 | 15 | 4 | 5 | 7 |
| 10 | 12 | 10 | 3 | 8 |

Fig. 2. Initial coding for machine $A$.


Fig. 3. Trial state assignment for machine $A$.
cannot be made with the constructed assignment, one may choose to experiment with a modification of the initial code before adding another state variable to the assignment.

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## Serial Adders with Overflow Correction

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Abstract-A method of implementing two single-bit adders is discussed. These adders can be used individually to realize the conventional functions of serial addition and serial multiplication on a pair of operands, or they can be cascaded to allow the serial addition of three operands for forming the product of complex numbers. In either case, the circuits will detect the occurrence of an overflow or the generation of the number minus one, and they will allow an addition to be rescaled by outputting the correct bits during the additional shifts, whether the addition overflowed or not.

Index Terms-Cascaded adders, overflow detection and correction, parallel processing, serial addition.

Manuscript received June 16, 1970; revised January 4, 1971.
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