choice of $\epsilon$. The usual procedure is to divide the range of integration into two parts, integrate outwards for a solution satisfying one boundary condition, integrate inwards for a solution satisfying the other boundary condition, match the solutions at an intermediate point and adjust $\epsilon$ so that the derivatives also agree [1], [2]. The inward integration may be avoided with the procedure described earlier. A convenient way of dividing the range is according to the sign of $f(r)$. For some $r$, $f(r)<0$ so that condition (ii) is not satisfied: the procedure described here is not always numerically stable when $f(r)<0$ [3]; in fact, for some values of $i,\left|d_{i}\right|<1$. Of a series of standard methods, the Numerov method,

$$
\begin{align*}
y_{n+1}=\left(\left(2+\frac{10}{12} h^{2} f_{n}\right) y_{n}-\right. & \left(1-\frac{h^{2}}{12} f_{n-1}\right) y_{n-1}  \tag{5}\\
& \left.+\frac{h^{2}}{12}\left(g_{n+1}+10 g_{n}+g_{n-1}\right)\right) /\left(1-\frac{h^{2}}{12} f_{n+1}\right),
\end{align*}
$$

was found to be most accurate in this case, for a given number of evaluations of $f$ for the outward integration. The procedure used successfully was to integrate outwards according to (5) until $f(r)>0$, then, with the last value computed as a boundary condition, to solve for the "tail" of the wave function by the method described here. The energy adjustment will be the same as before.

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# On the Inversion of Sparse Matrices 

By A. L. Dulmage and N. S. Mendelsohn

1. Introduction. There are a number of problems in applied mathematics involving many equations in many unknowns, but for which each equation involves only a small fraction of the unknowns. If such problems are linear or are approximated by linearization, one is involved with a matrix, a large proportion of whose entries are zero. To invert such a matrix $A$ it is sometimes advantageous to permute the rows and columns of $A$ yielding $P A Q$ where $P$ and $Q$ are permutation matrices. If

$$
P A Q=\left|\begin{array}{ccc}
A_{1} & & 0 \\
& A_{2} & \\
& & \\
* & & A_{r}
\end{array}\right|
$$

where $A_{1}, A_{2}, \cdots A_{r}$ are square matrices, the problem of inverting $P A Q$ is reduced

[^0]to inverting the blocks $A_{r}$ followed by matrix multiplication and addition. On permuting the rows and columns of $(P A Q)^{-1}$ one obtains the matrix $A^{-1}$.
F. Harary [4] gave a method based on the connectivity theory of directed graphs. The blocks $A_{i}$ turn out to be the matrix representatives of the strong components of an associated directed graph. Harary's method requires that $Q=P^{-1}$, a restriction which is quite unnecessary for matrix inversion. As a result, many matrices which reduce under independent permutations of rows and columns will not reduce if one insists that $Q=P^{-1}$. To remove this latter restriction, the authors have replaced a directed graph by a bipartite graph. The strong components of a directed graph become the irreducible components of a bipartite graph (see [1] for the definition of an irreducible component of a bipartite graph, [2] for the connection between strong components of directed graphs and irreducible components of bipartite graphs).
2. The Method of Reduction. Let $A$ be a square matrix of order $n$ with entries $a_{i j}$. Associated with $A$ is a bipartite graph $K_{A}$ with two sets of vertices $S=s_{1}$, $s_{2}, \cdots, s_{n}$ and $T=t_{1}, t_{2}, \cdots, t_{n}$. A pair $\left(s_{i}, t_{j}\right)$ is an edge of $K_{A}$ if and only if $a_{i j} \neq 0$ (one obtains Harary's directed graph if one identifies the vertices $s_{i}$ and $t_{i}$ ).

Suppose we can find matrices $P$ and $Q$ such that

$$
P A Q=\left|\begin{array}{ccc}
A_{1} & & 0 \\
& A_{2} & \\
& & \ddots \\
* & & A_{r}
\end{array}\right|
$$

and such that $A_{1}, A_{2}, \cdots, A_{r}$ cannot be further reduced by permutations of their rows and columns. Then the graphs corresponding to $A_{1}, A_{2}, \cdots, A_{r}$ are the irreducible components of $K_{A}$.

In [3], the authors give an algorithm for obtaining the irreducible components of a bipartite graph. This algorithm is easily programmable for machine computation and yields the permutations $P$ and $Q$. In the case where $A$ is a non-singular matrix the graph $K_{A}$ has no tails (see [3] for definition of a tail) and the algorithm described in [3] can be considerably simplified, since the steps needed to isolate and identify the tails can be omitted. An alternative procedure is the following. First, locate a nonzero term in the expansion of the determinant of $A$. This can readily be done using the algorithm of Marshall Hall [5] or Fulkerson and Ford [6]. Next, permute the rows of $A$ until the entries of this term occupy the main diagonal. Call this new matrix $A^{*}$. Finally, apply the method of F. Harary given in [4] to $A^{*}$.

We append an example:

$$
\text { Let } \quad A=\left|\begin{array}{ccccccc}
0 & 0 & 5 & 0 & 0 & 1 & 0 \\
0 & 2 & 0 & 0 & 0 & 0 & 4 \\
1 & 0 & 0 & 2 & 0 & 0 & 0 \\
0 & 4 & 0 & 3 & 0 & 0 & 0 \\
3 & 0 & 1 & 0 & 3 & 2 & 0 \\
3 & 2 & 0 & 0 & 2 & 0 & 0 \\
1 & 0 & 0 & 4 & 0 & 0 & 0
\end{array}\right|
$$

Note that the corresponding directed graph has the Hamiltonian circuit $1 \rightarrow 6 \rightarrow$ $5 \rightarrow 3 \rightarrow 4 \rightarrow 2 \rightarrow 7 \rightarrow 1$ and so is strongly connected. Hence, for no permutation $P$ does $P A P^{-1}$ reduce. Using the algorithm described in [3] one obtains the permutations $P=(1725643)$ and $Q=(1)(236754)$. Applying $P$ and $Q$ to the rows and columns of $A$, one obtains:

$$
P A Q=\left|\begin{array}{cc:ccc:cc}
1 & 2 & 0 & 0 & 0 & 0 & 0 \\
1 & 4 & 0 & 0 & 0 & 0 & 0 \\
\hdashline 0 & 3 & 4 & 0 & 0 & 0 & 0 \\
3 & 0 & 2 & 2 & 0 & 0 & 0 \\
0 & 0 & 2 & 0 & 4 & 0 & 0 \\
\hdashline 3 & 0 & 0 & 3 & 0 & 1 & 2 \\
0 & 0 & 0 & 0 & 0 & 5 & 1
\end{array}\right|
$$

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# Missing Data Correlation Computations 

By R. I. Jennrich

In correlation analysis or in any multivariate analysis based on the computation of a correlation or covariance matrix, the applied statistician often runs into the problem of missing data. To avoid complication in computing the correlation matrix, a complete observation vector is often discarded when only one or more of its components are missing. If a correlation matrix is computed by means of a standard electronic computer program, this procedure is often necessary. A large percentage of data may be thrown away when only a small percentage is missing. This note describes a modification in the standard computing scheme which eliminates this waste of data.

Let $x_{n 1}, x_{n 2}, \cdots, x_{n p}$ denote the $p$ components of the $n$th observation vector, $n=1,2, \cdots, N$. It is customary to add an $n+1$ st component to this vector which is identically equal to one. That is $x_{n, p+1}=1$. The cross product matrix

$$
a_{i j}=\sum_{n=1}^{N} x_{n i} x_{n j} \quad i=1, \cdots, p+1 ; \quad j=1, \cdots, p+1
$$

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