Algebraic elimination of slide surface constraints in implicit structural analysis

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SUMMARY

Slide surface and contact boundary conditions can be implemented via Lagrange multipliers in the algebraic equations in implicit structural analysis. This indefinite set of equations is difficult to solve by iterative methods and is often too large to be solved by direct methods. When there are \( m \) constraints and there exists a set of \( m \) variables where each variable is only involved in a single constraint, we advocate a direct elimination technique which leaves a sparse, positive definite system to solve by iterative methods. We prove that the amount of “fill-in” created by this process is independent of the size of the slide surfaces. In addition, the eigenvalues of the reduced matrix do not differ significantly from the eigenvalues of the unconstrained matrix. This method can be extended to the case where constrained surfaces intersect and leads to a graph theoretic approach for determining which variables can be eliminated efficiently for constraints with more general structure. Copyright © 2001 John Wiley & Sons, Ltd.

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1. INTRODUCTION

The numerical simulation of large structural deformations of materials has several applications such as the safety of high explosives and manufacturing problems in metal forming processes. During the heating of a high explosive, for instance, the explosive expands and deforms its containment vessel. Depending on the strength of the vessel, the explosive can achieve temperatures and pressures that cause it to reach a runaway state. In modeling this...
phenomenon, it is important to allow relative motion between the inside surface of the container and the interface defined by the explosive. Thus, the explosive and its container are modeled as two separate bodies with a *slide surface* boundary condition that allows the bodies to slide tangentially relative to each other.

Slide surface boundary conditions have been implemented in many explicit and implicit codes, e.g., [14, 5, 4]. In the implicit context, slide surface boundary conditions are implemented as constraints. The methods used to handle these constraints vary from Lagrange multipliers, e.g., [20], penalty methods, e.g., [14], or more recently, augmented Lagrange techniques [23, 17, 21]. In this paper we focus on the case of persistent contact and the use of Lagrange multipliers. The constraint equations are appended to the linear algebraic equations resulting from a Newton-Raphson linearization. The complete set of equations, unfortunately, is indefinite and is difficult to solve by iterative methods and is often too large to be solved by direct methods. The structure of the constraints, however, allows a subset of the equations to be eliminated directly, leaving a sparse, positive definite system to solve by iterative methods. This paper discusses the conditioning and sparseness of the reduced matrix and proposes a graph theoretic framework for determining which variables can be eliminated efficiently for constraints with more general structure.

In structural analysis, the “transformation method” (see, e.g., [8, 11]) and its variants [9, 1] are direct elimination methods. It seems that little attention has been given to partitioning the constrained degrees of freedom so that these methods are efficient, leading some authors to claim that these methods may be a “bottleneck” for large FE models [22]. Alternatively, Lagrange multiplier methods leading to KKT-type systems may be used, especially if direct solvers are available. If iterative solvers are used, for very large problems for example, the options are less attractive. Projection methods [12, 13, 22], Uzawa methods, e.g., [2, 10], and block preconditioners, e.g., [3, 16, 18, 15], may all suffer from slow convergence or high cost when many constraints are involved. When the constraints are structured such that direct elimination of some degrees of freedom is economical, then direct elimination as a preprocessing step can be effective before an iterative solution of the reduced system is carried out.

Section 2 of this paper briefly discusses the modeling of slide surfaces in the case of persistent contact. Section 3 derives an elimination technique for the Lagrange multiplier method which is equivalent to the transformation method. In Section 4, the sparseness and conditioning of the reduced matrix are discussed, as well as implementation options. Sample numerical results are given in Section 5. Section 6 shows how the elimination technique may be extended to problems with more general constraints. We close with some final remarks in Section 7.

## 2. SLIDE SURFACE MODELING

### 2.1. Slide surface constraints

A slide surface refers to the interface between two disjoint bodies that may come into contact in a simulation. Impact and separation of these bodies must be detected, and bodies may slide tangentially relative to each other, with or without friction. The slide surface boundary conditions are different in each case of impact, separation, and sliding. In this paper, we will only discuss modeling the simplest case, that of two bodies in contact that remain in contact, which may slide relative to each other during the simulation.
A slide surface constraint in this case is an “impenetrability” constraint preventing structural domains from overlapping. When it is known that two bodies are in contact, this constraint can be implemented by constraining one side of the slide surface, called the slave, to the other side, called the master. Since the master side then effectively defines the surface, it is typically chosen to be that side that moves less, is denser, more rigid, or more densely gridded. Interface nodes on the slave side are called “slave nodes,” and interface nodes on the master side are called “master nodes.” The surface defined by the master nodes is called the “master surface.”

In the case of persistent contact, the impenetrability constraint for a particular slave node constrains its displacement $u_s$ to lie on the master surface. Given a displacement $u_s^*$ that lies on the master surface near $u_s$, the impenetrability constraint states that

$$n^T(u_s - u_s^*) = 0$$

where $n$ is the normal to the master surface at $u_s^*$ (see Figure 1).

The parametric coordinates of $u_s^*$ are found by an iterative procedure that projects the current $u_s$ along the normal $n$ onto the closest element face on the master surface. The actual displacement of this point $u_s^*$ is then interpolated using the finite element shape functions. For linear hexahedral elements, the interpolation on the element face reduces to

$$u_s^* = \phi_1 u_1 + \phi_2 u_2 + \phi_3 u_3 + \phi_4 u_4$$

where $u_1$, $u_2$, $u_3$, and $u_4$ are the displacements of the four nodes defining the master face, and $\phi_1$, $\phi_2$, $\phi_3$, and $\phi_4$ are the corresponding shape function values for the parametric coordinates of $u_s^*$. Combining this with (1) yields

$$n^T(u_s - \phi_1 u_1 - \phi_2 u_2 - \phi_3 u_3 - \phi_4 u_4) = 0.$$ 

This constraint is nonlinear; the normal vector and the shape functions depend on both the master and the slave node displacements. We linearize this constraint by evaluating the normal and the shape functions using the current estimates for the displacements, while $u_s$, $u_1$, $u_2$, $u_3$, and $u_4$ represent the new displacements. We are currently experimenting with other linearizations of this constraint which may perform better in cases of large relative sliding.
2.2. Structure of the constraints

In matrix form, the set of constraints for all slave nodes is \( G^T u = 0 \), where \( u \) is the vector of nodal displacements. Assuming that the nodes are numbered such that all slave nodes follow all master nodes which follow all other nodes, the matrix \( G^T \) has a structure of the form

\[
\begin{bmatrix}
0 & \ldots & 0 & -\phi_{11}n_1^T & -\phi_{12}n_1^T & -\phi_{13}n_1^T & -\phi_{14}n_1^T \\
\vdots & \ddots & \vdots & \vdots & \vdots & \vdots & \vdots \\
0 & \ldots & 0 & -\phi_{22}n_2^T & -\phi_{23}n_2^T & -\phi_{24}n_2^T & -\phi_{25}n_2^T \\
\end{bmatrix}
\begin{bmatrix}
\mathbf{n}_1^T \\
\mathbf{n}_2^T \\
\vdots \\
\end{bmatrix}
\]

(2)

where \( n_i^T = (n_{ix}^i, n_{iy}^i, n_{iz}^i) \) is the unit normal for the \( i \)th slave node, and \( \phi_{ij} \) are shape function values for the four master nodes (indexed by \( j \)) associated with the same slave node. We abbreviate the above structure by \( [0, -W^T W, W^T] \) where \( W \) is a matrix containing the positive shape function values \( \phi_{ij} \). Since no node depends on any slave nodes in a slide surface constraint, the matrix \( N^T \) contains at most one nonzero in each column. The number of constraints is equal to the number of slave nodes on the slide surface and is typically much smaller than the total number of nodes.

3. ALGEBRAIC ELIMINATION FOR CONSTRAINED PROBLEMS

The discrete and linear constraints \( G^T u = 0 \) may be incorporated into the finite element equations \( Au = f \) via the transformation method or the Lagrange multiplier method. The latter leads to a system that is indefinite and generally difficult to solve by iterative methods. However, like the transformation method, solving this system can be reduced to solving a sparse, positive definite system. The economy of this transformation depends on the constraints having the specific structure described below.

3.1. Transformation method

In the transformation method [8, 11], the variables are partitioned into two sets, called independent and dependent. The number of dependent variables equals the number of constraint equations, \( m \). Suppose such a partitioning has been made, and let \( u_1 \) denote the independent variables and \( u_2 \) denote the dependent variables. The system \( Au = f \) may now be partitioned symmetrically into block form as

\[
\begin{bmatrix}
A_{11} & A_{12} \\
A_{21} & A_{22}
\end{bmatrix}
\begin{bmatrix}
u_1 \\
u_2
\end{bmatrix}
= \begin{bmatrix}
f_1 \\
f_2
\end{bmatrix}
\]

while \( G^T u = 0 \) can be partitioned as

\[
\begin{bmatrix}
B^T & D^T
\end{bmatrix}
\begin{bmatrix}
u_1 \\
u_2
\end{bmatrix}
= 0.
\]

The transformation matrix \( T \) is defined so that \( u = Tu_1 \), hence

\[
T = \begin{bmatrix}
I \\
-D^{-T} B^T
\end{bmatrix}.
\]

The transformed system to be solved is

\[
T^T A Tu_1 = T^T f
\]
the solution of which may be substituted into \( u = T u_1 \) to recover the full solution. The transformed coefficient matrix is

\[
T^T A T = A_{11} - A_{12} D^{-T} B T - B D^{-1} A_{21} + B D^{-1} A_{22} D^{-T} B^T
\]

which is symmetric positive definite by construction, provided \( A \) is symmetric positive definite.

We now return to the problem of partitioning the variables into independent and dependent sets such that (3) is economical to form. We define a purely dependent variable to be a variable that is only involved in one constraint. If all the dependent variables \( u_2 \) are purely dependent, then \( D \) is composed of one nonzero per column and one nonzero per row (it is a reordered diagonal matrix) and its inverse is sparse. The transformed matrix (3) would then be sparse and economical to construct and factorize. In many problems it is possible to choose \( m \) purely dependent variables. This is often the case in structural analysis since slave variables are often defined so that they are purely dependent.

For the constraint matrix (2), each column of \( N^T \) has at most one nonzero entry, and since each row of \( N^T \) represents a unit vector, each row contains at least one nonzero entry. Choosing the \( m \) purely dependent variables amounts to choosing one of the components of \( n^T_i = (n^x_i, n^y_i, n^z_i) \) for each constraint \( i \). Thus we can reorder and partition \( N^T \) by columns into \( \begin{bmatrix} N^T_1 \ N^T_2 \end{bmatrix} \) such that \( N^T_2 \) is a diagonal matrix.

3.2. Lagrange multiplier method

When the variables are partitioned into slave nodes, master nodes, and interior (all other) nodes, the Lagrange multiplier method leads to an equation with the structure

\[
\begin{bmatrix}
A_i & A_{im} & A_{is} \\
A_{mi} & A_m & -W^T N \\
-A_s & N^T & 0
\end{bmatrix}
\begin{bmatrix}
u_i \\
u_m \\
u_s
\end{bmatrix} =
\begin{bmatrix}
f_i \\
f_m \\
\lambda
\end{bmatrix}
\]

(4)

where \( \lambda \) is the vector of Lagrange multipliers and can be interpreted as the normal force at each slave node necessary to conserve momentum. There is no interaction in the stiffness matrix \( A \) between the slave nodes and the master nodes.

A technique equivalent to the transformation method may be derived for matrices in this form. As before, we partition the constrained variables into independent and dependent sets. For our slide surface constraints, we partition \( G^T \) into \( \begin{bmatrix} 0, -N^T W, N^T_1, N^T_2 \end{bmatrix} \) where \( N^T_2 \) corresponds to the dependent set. We can now further partition (4) as

\[
\begin{bmatrix}
X & X & X \\
X & X & X \\
X & X & A_{22} \\
-X^T W & N^T_1 & N^T_2
\end{bmatrix}
\begin{bmatrix}
u_i \\
u_m \\
u_{s_1} \\
\lambda
\end{bmatrix} =
\begin{bmatrix}
f_i \\
f_m \\
f_{s_1} \\
0
\end{bmatrix}
\]

(5)

where \( X \) represents a nonzero block in the matrix and \( u_{s_1} \) and \( u_{s_2} \) are slave variables. We will use the following simplified notation for (5):

\[
\begin{bmatrix}
A_{11} & A_{12} & B \\
A_{21} & A_{22} & D \\
B^T & D^T & 0
\end{bmatrix}
\begin{bmatrix}
u_1 \\
u_2 \\
\lambda
\end{bmatrix} =
\begin{bmatrix}
f_1 \\
f_2 \\
0
\end{bmatrix}
\]

(6)
The Schur complement of the matrix in (6) with respect to \[
\begin{bmatrix}
A_{22} & D
\end{bmatrix}
\] is
\[
S = A_{11} - \begin{bmatrix}
A_{12} & B
\end{bmatrix} \begin{bmatrix}
A_{22} & D \\
D^T & 0
\end{bmatrix}^{-1} \begin{bmatrix}
A_{21} \\
B^T
\end{bmatrix}.
\] (7)

Again, if there are \(m\) purely dependent variables that can be chosen for the dependent set, then it is possible to choose \(D\) such that it is diagonal. In this case, the inverse
\[
\begin{bmatrix}
A_{22} & D \\
D^T & 0
\end{bmatrix}^{-1} = \begin{bmatrix}
0 & D^{-T} \\
D^{-1} & -D^{-1}A_{22}D^{-T}
\end{bmatrix}
\] (8)
is sparse, which generally leads to a sparse Schur complement. The Schur complement is positive definite since it is equivalent to the transformed matrix (3).

The block system (6) can be solved by solving the reduced system
\[
Su_1 = f_1 - BD^{-1}f_2
\] (9)
and substituting \(u_1\) into
\[
u_2 = -D^{-T}B^Tu_1
\] (10)
\[
\lambda = D^{-1}(f_2 - A_{21}u_1) + D^{-1}A_{22}D^{-T}B^Tu_1.
\] (11)

4. PROPERTIES OF THE REDUCED MATRIX

4.1. Eigenvalues of \(S\)

The convergence rate of the conjugate gradient method for solving systems with the reduced matrix \(S\) depends on the eigenvalues of \(S\). If \(A\), the stiffness matrix without the constraints, is symmetric positive definite, the next lemma bounds the eigenvalues of \(S\) away from the origin.

**Lemma 4.1.** Define \(A = \begin{bmatrix}
A_{11} & A_{12} \\
A_{21} & A_{22}
\end{bmatrix}\) and \(S\) as in equation (7). If \(A\) is symmetric, the smallest eigenvalue of \(S\), \(\lambda_{\min}(S)\), is larger than the smallest eigenvalue of \(A\), \(\lambda_{\min}(A)\).

**Proof.** Choose \(x_1 \neq 0\) such that
\[
\frac{x_1^T S x_1}{x_1^T x_1} = \lambda_{\min}(S)
\]
and let \(x = \begin{bmatrix} x_1 \\ x_2 \end{bmatrix}\), \(x_2 = -D^{-T}B^Tx_1\). Then,
\[
\frac{x_1^T S x_1}{x_1^T x_1} \geq \frac{x_1^T S x_1}{x_1^T x_1 + x_2^T x_2} = \frac{x^T A x}{x^T x} \geq \lambda_{\min}(A).
\]

\(\square\)

The behavior of the eigenvalues can be determined by applying well-known theorems in eigenvalue sensitivity to the reduced matrix in the form (7). Consider the symmetric matrix
\[
A = \begin{bmatrix}
A_{11} & B \\
B^T & W
\end{bmatrix}, \quad \text{with} \quad W = \begin{bmatrix}
a \\ d
\end{bmatrix}, \quad a > 0, \quad d \neq 0
\]
and the pairwise elimination of a single constraint and its corresponding slave equation to get the Schur complement \( S \). (Eliminating \( m \) constraints and their corresponding slave equations is simply a sequence of these pairwise eliminations.)

Define \( E \equiv -BW^{-1}B^T \) so that \( S = A_{11} + E \). The eigenvalues of \(-W^{-1}\) are

\[
\frac{a \pm \sqrt{a^2 + 4d^2}}{2d^2}
\]

which shows that one eigenvalue is positive and the other is negative. From the Sylvester law of inertia, assuming \( \bar{B} \) has full rank, \( E \) has one positive eigenvalue and one negative eigenvalue; the other eigenvalues are zero.

For a matrix \( S \), let \( \lambda_k(S) \) denote the \( k \)th largest eigenvalue and let \( n \) denote the dimension of \( S \). Two corollaries of the Courant-Fischer minimax theorem apply. First, the extremal eigenvalues of \( A_{11} \) are bounded by the extremal eigenvalues of \( A \) (interlacing property). Second, we have (see [24])

\[
\begin{align*}
\lambda_{r+s-1}(S) & \leq \lambda_r(A_{11}) + \lambda_s(E), \quad r + s - 1 \leq n \\
\text{and} \quad \lambda_{r-s+1}(S) & \geq \lambda_r(A_{11}) + \lambda_{n-r+1}(E), \quad s \leq r \leq n.
\end{align*}
\]

In particular, if the positive eigenvalue of \( E \) is bounded, then the largest eigenvalue of \( S \) is bounded. In addition, if \( E \) has a very large positive eigenvalue, the bound on the largest eigenvalue of \( S \) increases by the size of this eigenvalue, but since \( E \) has many zero eigenvalues, the other eigenvalues of \( S \) are still bounded by the eigenvalues of \( A_{11} \),

\[
\begin{align*}
\lambda_{r+1}(S) & \leq \lambda_r(A_{11}), \quad 1 \leq r \leq n - 1 \\
\lambda_{r+1}(S) & \geq \lambda_{r+2}(A_{11}), \quad 0 \leq r \leq n - 2.
\end{align*}
\]

Small eigenvalues of \( E \) can generally be avoided by choosing to eliminate a slave variable that corresponds to large magnitude values of \( d \), and, less importantly, small values of \( a \). (This assumes that the nonsquare matrix \( \bar{B} \) is not too poorly conditioned.) For slide surface constraints, small magnitude values of \( d \) can easily be avoided since \( |d| \geq 1/\sqrt{3} \) can always be chosen (as a component in a 3-D unit normal vector).

Instead of being concerned by the eigenvalues of \( E \), we can guarantee that the spectrum of \( S \) does not differ too much from the spectrum of \( A_{11} \) by bounding the size of the entries in \( E \), or more precisely, the Frobenius norm of \( E \). This is equivalent to bounding the sum of the squares of the eigenvalues of \( E \) and is the result of the Wielandt-Hoffmann theorem,

\[
\sum_{i=1}^{n} (\lambda_i(S) - \lambda_i(A_{11}))^2 \leq \|E\|_F^2.
\]

Again, the norm of \( E \) can be roughly controlled by not choosing excessively small magnitude values of \( d \).

It is possible in some cases to choose variables to eliminate such that the conditioning of \( S \) is better than the conditioning of \( A_{11} \). This is very difficult to do in general, however.

In conclusion, we try to control the spectrum of \( S \) by guaranteeing that it does not differ too much from the spectrum of \( A_{11} \). Very small magnitudes of \( d \) can seriously harm convergence and can easily be avoided by simply choosing larger components in \( n_i^T = (n_i^x, n_i^y, n_i^z) \) for each constraint row \( i \) when partitioning \( N^T \). With this strategy, we find in practice that solving with \( S \) is only slightly more difficult than solving with \( A \). Section 5 will show a typical example of this.
4.2. Sparseness of $S$

In this section we show that the number of nonzeros in each row of $S$ is bounded, and is not dependent on the number of slave variables or constraints. Thus the algebraic elimination method can be used for very large problems. In the following, we will use the graph theoretic representation of a sparse matrix. To analyze the sparsity pattern of $S$, “fill-paths” and the graph interpretation of Gaussian elimination [19] give an over-estimate of the fill-in that occurs, since they assume that the inverse (8) is dense. Hence, we proceed as follows.

The variables in $u$ can be partitioned into interior, master, slave, uneliminated slave, and constraint (Lagrange multiplier), and these sets may be represented by the symbols $i$, $m$, $s$, $t$, and $c$, respectively. Uneliminated slave variables are those variables at the same grid point as an eliminated slave variable. We then label the relevant blocks in (5) as follows:

$$
\begin{bmatrix}
X & X & X & A_{is} & A_{mc} \\
X & X & A_{is} & A_{ts} & A_{tc} \\
A_{si} & A_{st} & A_{ss} & A_{sc} \\
A_{cm} & A_{ct} & A_{cs} & 0
\end{bmatrix}.
$$

(12)

The sparsity pattern of a matrix may be described by its graph. For a symmetric matrix $A$ of order $n$, its graph $G(A)$ is composed of the vertices $\{1, \ldots, n\}$ and the edges $\{(i,j), i \neq j \text{ and } A_{ij} \neq 0\}$. When the context is clear, we will not distinguish between vertices and variables, and between edges and nonzeros. The graph of a matrix is often related to its discretization grid. An example graph (for a matrix without uneliminated slave variables, i.e., a scalar problem) is shown in Figure 2.
Figure 3. Graph of the reduced matrix after the slave and constraint vertices have been eliminated. Dotted lines indicate the fill-in edges. Note that there is no fill-in edge between nodes 9 and 11 as would be predicted by the fill-path theorem for Gaussian elimination.

For the matrix (12), the “fill-in matrix,” $S - A_{11}$, after eliminating the $s$ and $c$ variables is

$$
\begin{bmatrix}
0 & A_{is}A_{sc}^{-1}A_{cm} & A_{is}A_{sc}^{-1}A_{ct} \\
A_{mc}A_{cs}^{-1}A_{si} & -A_{mc}A_{cs}^{-1}A_{ss}A_{sc}^{-1}A_{cm} & A_{mc}A_{cs}^{-1}A_{st} \\
A_{tc}A_{cs}^{-1}A_{si} & -A_{tc}A_{cs}^{-1}A_{ss}A_{sc}^{-1}A_{cm} & A_{tc}A_{cs}^{-1}A_{st}
\end{bmatrix}.
$$

In the graph of this matrix, we first note that there is no fill-in between any interior vertices. The $(1,2)$ block $A_{is}A_{sc}^{-1}A_{cm}$ and the $(2,1)$ block $A_{mc}A_{cs}^{-1}A_{si}$ correspond to fill-in between an interior vertex and a master vertex when there is path between these vertices through a slave vertex and a constraint vertex in $G(A)$. These paths may be abbreviated $(i,s,c,m)$ and $(m,c,s,i)$. The $(2,2)$ block $-A_{mc}A_{cs}^{-1}A_{ss}A_{sc}^{-1}A_{cm}$ corresponds to fill-in between master vertices that are connected via a path $(m,c,s,s,c,m)$.

The fill-in caused by the other blocks can be interpreted similarly. None of these blocks cause fill-in between vertices that are more than a small fixed number of path lengths apart. Thus the fill-in in a row of $S$ is bounded.

Figure 3 shows the graph of $S$ for the problem of Figure 2. The fill-in is shown with dotted edges. The edge $(7,9)$ is caused by a path of the type $(m,c,s,s,c,m)$. The fill-in generally involves vertices on or adjacent to the slide surface.

Counting paths in the graph of (13) gives the number of fill-ins that will occur for a given slide surface geometry.

### 4.3. Alternative implementations

To solve a system with $S$ by iterative methods, it is not necessary to form $S$ but only to apply its action to a vector. Even if $S$ does need to be formed, for example for use by a preconditioner,
this cost is low since $S$ is sparse as shown above. Forming the Schur complement directly via (3) is stable in the sense that the process will not break down. The size of the largest entries in $S$ is controlled insofar as the smallest entries in $D$ are controlled. As mentioned, for slide surface constraints we can choose components in $D$ as the largest component of a unit vector in 3-D, and these components are at least $1/\sqrt{3}$ in magnitude.

Besides forming $S$ directly, it is conceivable to adapt a direct solver to construct $S$ in order to save programming effort. If the matrix in (6) is symmetrically reordered such that the first $2m$ equations are the individual slave equations to be eliminated followed immediately by their corresponding constraint equations (i.e., equations corresponding to slaves and constraints are interlaced), then $S$ is the Schur complement that remains after the first $2m$ variables are eliminated via either an LU factorization or symmetric indefinite factorization with 2-by-2 blocks, both without pivoting. This ordering reduces intermediate fill-in in the LU factorization.

5. NUMERICAL TESTS

An octant of three concentric spherical shells is shown in Figure 4. The first and third shells are composed of steel, and the second shell is composed of lucite. Slide surface boundary conditions are used between the steel and lucite shells. This test problem has 11400 elements and 13832 nodes, with 3 displacement variables at each node. The same problem without modeling the slide surfaces has 12844 nodes. In this example, the lucite is given an initial internal energy of 1 kbar-cc, causing it to expand and slide relative to the steel shells. The sliding is small, so this is a test of the method in the case of linear persistent contact.

The algebraic elimination technique was applied to the linear systems that arise and the resulting reduced systems were solved with the preconditioned conjugate gradient method with factorized sparse approximate inverse preconditioning [6, 7]. A zero initial guess was used, and the residual norm was reduced by six orders of magnitude. Four processors of an IBM SP computer were used in the computation, with interprocessor communication implemented.
with MPI. In the following table, we compare a few performance metrics between problems with and without slide surface constraints.

Table I. Comparison between problems with and without slide surfaces.

<table>
<thead>
<tr>
<th></th>
<th>With slide surfaces</th>
<th>Without slide surfaces</th>
</tr>
</thead>
<tbody>
<tr>
<td>Nonzeros in $S$ (with slides) or $A$ (without)</td>
<td>1.64 million</td>
<td>1.37 million</td>
</tr>
<tr>
<td>Nonzeros in preconditioner</td>
<td>0.95 million</td>
<td>0.91 million</td>
</tr>
<tr>
<td>Algebraic elimination time</td>
<td>0.9 seconds</td>
<td>N/A</td>
</tr>
<tr>
<td>Preconditioner setup time</td>
<td>9.3 seconds</td>
<td>6.0 seconds</td>
</tr>
<tr>
<td>Iteration count</td>
<td>646</td>
<td>646</td>
</tr>
<tr>
<td>Total solution time</td>
<td>33.6 seconds</td>
<td>26.7 seconds</td>
</tr>
</tbody>
</table>

The table shows that the reduced matrix $S$ contains only slightly more nonzeros than the matrix $A$ for a problem that does not model slide surfaces. The iteration counts are also similar (in this case they happen to be the same). In addition, the time for the elimination step is only a small fraction of the overall solution time.

6. EXTENSION TO MORE GENERAL CONSTRAINTS

6.1. Intersecting constrained surfaces

For general constraints, it is not always possible to find a partitioning of the constrained variables into independent and dependent sets such that the dependent variables are purely dependent. Algebraically, this means that it may be impossible to reorder the columns and partition $G^T$ into $[B^T, D^T]$ such that $D$ is diagonal. However, the algebraic elimination procedure can still be economical if the inverse of $D$ is sparse, e.g., when $D$ is a block diagonal matrix.

A complete set of purely dependent variables cannot be found, for example, when two slide surfaces intersect at a T. Figure 5 shows this case, which is actually treated as three slide surfaces. Master and slave sides are chosen for each slide surface, and these are marked in the figure. Arrows in the figure indicate that a slave node is dependent on a master node. Node 2 is a slave node for one slide surface, but is a master node for another. Also, node 5 is a slave node for two slide surfaces. Nodes 2 and 5 are not purely dependent nodes.

For intersecting slide surfaces, however, it is possible to partition the variables such that $D$ is block diagonal. Consider the following procedure for selecting the dependent variables, which corresponds to choosing columns of $G^T$ to form $D$. In this discussion, we will not differentiate between variables and nodes, i.e., we consider a scalar problem.

We start by choosing all slave nodes to be in the dependent set. If a slave node is not purely dependent, then it falls under one of these two cases:

1. the node is a slave in one constraint and a master in one or more other constraints,
2. the node is a slave node in more than one constraint.

In the first case, if a slave node is also a master in $k$ constraints, then $D$ will contain a block
of size \((k + 1)\times(k + 1)\) since the column of \(G^T\) corresponding to the slave node has \(k + 1\) nonzeros. Nothing needs to be done in this case.

In the second case, there are more constraints than slave nodes and one or more master nodes must also be selected for the dependent set. Let \(s_1\) be a slave node that participates in constraints \(c_1\) and \(c_2\). A master node must be chosen from the master nodes in \(c_1\) or \(c_2\) so that \(D\) will be nonsingular. If the master node is a node that participates in \(k\) constraints, then \(D\) will contain block of size \(k\times k\); master nodes with small \(k\) should be preferred. If \(s_1\) is a slave node in more than two constraints, then two or more master nodes need to be selected simultaneously. Here, master nodes should be selected so that the total number of new constraints involved is minimized.

In summary, the dependent set should be composed of all the slave nodes, and selected master nodes as described in the second case. The blocks formed in \(D\) may overlap so that the actual blocks may be larger than the sizes mentioned. This procedure is effective, but because of possible overlapping blocks, it is not optimal in the sense that the smallest blocks are found.

### 6.2. General constraints

In the case where the constraints are even more general, a block diagonal \(D\) can often still be found. In this section, we develop a graph theoretic framework for partitioning \(G^T\) so that \(D\) is block diagonal with small blocks.

Define the binary matrix \(Q^T\) such that \(\{Q^T\}_{ij} = 1\) if and only if constraint \(i\) involves node \(j\) (which is either a master or slave node). The matrix \(Q^T\) contains \(m\) rows and \(n\) columns, where \(m\) is the number of constraints and \(n\) is the total number of slave and master nodes. For the geometry in Figure 5, we have the matrix

\[
Q^T = \begin{bmatrix}
1 & 1 & . & 1 & . & . & . & . & . & . \\
1 & 1 & . & . & 1 & . & . & . & . & . \\
. & . & . & . & 1 & . & . & 1 & 1 & . \\
. & . & . & . & . & 1 & . & 1 & 1 & . \\
. & 1 & . & . & . & 1 & 1 & . & . & . \\
. & . & . & 1 & . & . & 1 & . & . & . \\
. & . & . & . & . & . & . & 1 & 1 & . \\
. & . & . & . & . & 1 & . & . & 1 & . \\
. & . & . & . & . & . & 1 & . & . & 1 \\
. & . & . & . & . & . & . & 1 & . & 1 \\
\end{bmatrix}
\]  

(14)
The matrix $QQ^T$ contains a nonzero at $(i,j)$ if and only if nodes $i$ and $j$ are involved in the same constraint. The graph of $QQ^T$, denoted by $G(QQ^T)$, will be helpful to understand the partitioning problem at hand. For $Q^T$ given by (14), $G(QQ^T)$ is shown in Figure 6.

A clique is a subgraph such that every vertex in the subgraph has an edge to every other vertex in the subgraph. A maximal clique is a clique such that no vertex can be added to the clique to form a larger clique. The graph $G(QQ^T)$ contains $m$ maximal cliques, each corresponding to a constraint, and the vertices of each clique correspond to the nodes involved in the constraint. The partitioning problem is to select $m$ nodes (or equivalently, vertices) such that $D$ is block diagonal with small blocks.

**Fact 6.1.** If an independent set of $m$ vertices exists in $G(QQ^T)$, then $Q^T$ can be partitioned into $[Q_1^T, Q_2^T]$ such that $Q_2^T$ is a diagonal matrix.

An independent set of size $m$ may not exist. Further, finding this independent set is equivalent to the maximum independent set problem, which is NP-hard.

**Fact 6.2.** A vertex from each maximal clique in $G(QQ^T)$ must be selected, otherwise $D$ will be singular.

**Fact 6.3.** A connected set of $k$ selected vertices in $G(QQ^T)$ will form a block of size exactly $k$-by-$k$ in $D$.

**Fact 6.4.** Let $C_j$ denote the number of maximal cliques in $G(QQ^T)$ that vertex $j$ belongs to. Selecting this vertex for the dependent set will form a block of size at least $C_j$-by-$C_j$ in $D$.

The partitioning problem may be solved, for example, by algorithms that try to select vertices with small $C_j$. For example, $m$ vertices with the smallest $C_j$ may be selected to form the dependent set. Figure 7 is a possible solution with this strategy for the graph of Figure 6. Three vertices with $C_j = 1$ are selected (solid circles) and three vertices with $C_j = 2$ are selected (gray circles). For comparison, Figure 8 shows a situation where all the vertices corresponding to slave nodes (solid circles) are selected first, as suggested in the previous subsection. Here, the solution is worse because $D$ will contain a large 4-by-4 block. Selecting vertices with small $C_j$ is a poor algorithm for many other problems, however. We intend to further study the partitioning problem using this graph theoretic framework in a future paper.
6.3. Dense rows in $G^T$

Dense rows in $G^T$ occur when constraints involve every variable of the problem, and make the algebraic elimination technique very costly. As an alternative, these rows may be omitted from the elimination. Consider the partitioning of the matrix in (6) into

\[
\begin{bmatrix}
\tilde{A} & \tilde{B} \\
\tilde{B}^T & 0
\end{bmatrix}
\]

where $\tilde{B}^T$ represents dense or relatively full rows, and $\tilde{A}$ represents the remainder of the problem with sparse constraints. The algebraic elimination technique is only applied to $\tilde{A}$. If the number of rows in $\tilde{B}^T$ is small, the reduced system based on the matrix $\tilde{B}^T \tilde{A}^{-1} \tilde{B}$ may be solved iteratively. The reduced matrix is not formed, and an iterative method for indefinite systems is required. Each iteration involves a solve with $\tilde{A}$ which in turn involves a solve with a positive definite matrix. If the number of rows in $\tilde{B}^T$ is very small as in many applications, then very few iterations may be required.
7. CONCLUDING REMARKS

We have been regularly using the algebraic elimination technique described in this paper for symmetric problems, and have not encountered any difficulties with very poorly conditioned reduced matrices. The extension to more general constraints in Section 6 may be very valuable, and in future research we intend to develop and test algorithms based on these ideas.

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