# Determination of the $D^{1 / 2}$-Norm of the SOR Iterative Matrix for the Unsymmetric Case 

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#### Abstract

This paper is concerned with the determination of the Jordan canonical form and $D^{1 / 2}$-norm of the SOR iterative matrix derived from the coefficient matrix $A$ having the form $$
A=\left(\begin{array}{cc} D_{1} & -H \\ H^{T} & D_{2} \end{array}\right)
$$ with $D_{1}$ and $D_{2}$ symmetric and positive definite. The theoretical results show that the Jordan form is not diagonal, but has only $q$ principal vectors of grade 2 and that the $D^{1 / 2}$-norm of $\mathscr{L}_{\omega_{b}}$ ( $\omega_{b}$, the optimum parameter) is less than unity if and only if $\bar{\mu}=\rho(B)$, the spectral radius of the associated Jacobi iterative matrix, is less than unity. Here $q$ is the multiplicity of the eigenvalue $i \bar{\mu}$ of $B$.


1. Introduction. For the iterative solution of the linear system of equations,

$$
\begin{equation*}
A x=b \tag{1.1}
\end{equation*}
$$

the Jacobi and Gauss-Seidel methods are well known. They are very simple from a computational point of view since only matrix-vector multiplications and linear combinations of vectors are needed. This is also valid for the modification called "Successive Overrelaxation" or "SOR" method, where a relaxation factor is introduced for accelerating the convergence. Let

$$
\begin{equation*}
A=D-A_{L}-A_{U} \tag{1.2}
\end{equation*}
$$

where $D$ is the block diagonal part of $A,-A_{L}$ and $-A_{U}$ are the remaining strictly lower and upper triangular parts of $A$; then, if $D$ is nonsingular, the SOR method is given by

$$
\begin{equation*}
x_{k+1}=\mathscr{L}_{\omega} x_{k}+\omega(I-\omega L)^{-1} D^{-1} b, \quad k \geq 0 \tag{1.3}
\end{equation*}
$$

Here, $x_{0}$ is an initial vector, $\mathscr{L}_{\omega}$ the iterative matrix given by

$$
\begin{equation*}
\mathscr{L}_{\omega}=(I-\omega L)^{-1}[(1-\omega) I+\omega U] \tag{1.4}
\end{equation*}
$$

and

$$
\begin{equation*}
L=D^{-1} A_{L}, \quad U=D^{-1} A_{U} \tag{1.5}
\end{equation*}
$$

Now (1.3) converges if and only if the spectral radius of $\mathscr{L}_{\omega}$ is less than unity, and the asymptotic rate of convergence is given by

$$
\begin{equation*}
R_{\infty}\left(\mathscr{L}_{\omega}\right)=-\log \left(\rho\left(\mathscr{L}_{\omega}\right)\right) \tag{1.6}
\end{equation*}
$$

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The SOR method has been extensively studied for a symmetric positive definite matrix $A$ (see, e.g., Varga [5] and Young [7]). For a positive definite and consistently ordered matrix $A$, from [5] and [7] we have:

S1. $\rho\left(\mathscr{L}_{\omega}\right)<1 \Leftrightarrow \bar{\mu}<1$ and $0<\omega<2$.
S2.

$$
\rho\left(\mathscr{L}_{\omega}\right)= \begin{cases}\left\{\omega \bar{\mu}+\left[\omega^{2} \bar{\mu}^{2}-4(\omega-1)^{1 / 2}\right]^{1 / 2}\right\}^{2} / 4 & \text { if } 0<\omega<\omega_{b}^{\prime}, \\ \omega-1 & \text { if } \omega_{b}^{\prime} \leq \omega<2 .\end{cases}
$$

S3. $\rho\left(\mathscr{L}_{\omega_{b}^{\prime}}\right)<\rho\left(\mathscr{L}_{\omega}\right)$ if $\omega \neq \omega_{b}^{\prime}$.
Here,

$$
\bar{\mu}=\rho(B)=\rho(L+U), \quad \omega_{b}^{\prime}=2 /\left[1+\left(1-\bar{\mu}^{2}\right)^{1 / 2}\right] .
$$

Young [6] has shown that if $A$ is consistently ordered and the eigenvalues of $B$ are real and less than unity in modulus, then the Jordan canonical form of $\mathscr{L}_{\omega_{b}^{\prime}}$ is not diagonal. Therefore, in this case, the SOR method converges slower than expected based on the spectral radius $\rho\left(\mathscr{L}_{\omega_{b}^{\prime}}\right)$. When $A$ is symmetric and positive definite, and when $A$ has the form

$$
A=\left(\begin{array}{cc}
D_{1} & -H  \tag{1.7}\\
-K & D_{2}
\end{array}\right)
$$

Young [7, Chapter 7] determined the $D^{1 / 2}$-norm and $A^{1 / 2}$-norm of $\mathscr{L}_{\omega_{b}^{\prime}}$ (the spectral norms of $D^{1 / 2} \mathscr{L}_{\omega_{b}^{\prime}} D^{-1 / 2}$ and $A^{1 / 2} \mathscr{L}_{\omega_{b}^{\prime}} A^{-1 / 2}$, respectively) and pointed out that the $D^{1 / 2}$-norm of $\mathscr{L}_{\omega_{b}^{\prime}}$ is greater than unity in general. Moreover, $\left\|\mathscr{L}_{\omega_{b}^{\prime}}^{m}\right\|$ (the spectral norm of $\left.\mathscr{L}_{\omega_{b}^{\prime}}^{m}\right)$ behaves much like $\left\|\mathscr{L}_{\omega_{b}^{\prime}}^{m}\right\|_{D^{1 / 2}}$. However, for large $m,\left\|\mathscr{L}_{\omega_{b}^{\prime}}^{m}\right\|_{D^{1 / 2}}<1$, and eventually $\left\|\mathscr{L}_{\omega_{b}^{\prime}}^{m}\right\|_{D^{1 / 2}}$ tends to zero, though considerably more slowly than $\rho\left(\mathscr{L}_{\omega_{b}^{\prime}}\right)$.

For the matrix $A$ in (1.1) the unsymmetric case is by far not as common as the symmetric one, but nevertheless, unsymmetric matrices appear, e.g., in the numerical solution of the biharmonic equation [1] and the computation of cubic splines, [3] and [4, Chapter 3]. If the matrix $A(1.2)$ is consistently ordered and $B$, given by

$$
\begin{equation*}
B=L+U \tag{1.8}
\end{equation*}
$$

is similar to a skew-symmetric matrix and has either zero eigenvalues or purely imaginary eigenvalues, then from [1], [3], and [4], or the theory of Young [7], we have:

US1. $\rho\left(\mathscr{L}_{\omega}\right)<1 \Leftrightarrow 0<\omega<2 /(1+\bar{\mu})$.
US2.

$$
\rho\left(\mathscr{L}_{\omega}\right)= \begin{cases}1-\omega & \text { if } 0<\omega \leq \omega_{b}, \\ {\left[\frac{\bar{\mu} \omega+\sqrt{\omega^{2} \bar{\mu}^{2}+4(\omega-1)}}{2}\right]^{2}} & \text { if } \omega_{b}<\omega<\frac{2}{1+\bar{\mu}} .\end{cases}
$$

US3. $\rho\left(\mathscr{L}_{\omega_{b}}\right)<\rho\left(\mathscr{L}_{\omega}\right)$ if $\omega \neq \omega_{b}$.
Here,

$$
\begin{equation*}
\bar{\mu}=\rho(B), \quad \omega_{b}=2 /\left(1+\sqrt{1+\bar{\mu}^{2}}\right) \tag{1.9}
\end{equation*}
$$

Notice that in this case we can always choose the relaxation factor $\omega$ such that $\rho\left(\mathscr{L}_{\omega}\right)<1$, no matter how large $\bar{\mu}$ is. This is very different from the symmetric
case. Another difference between the two cases is that the optimum factor $\omega_{b}$ for the unsymmetric case is less than unity and the optimum factor $\omega_{b}^{\prime}$ for the symmetric case is greater than unity. It is also important to note that overestimating $\omega_{b}^{\prime}$ is better than an underestimation, but for $\omega_{b}$ an underestimate is better than overestimating.

However, to our knowledge, the Jordan canonical form and $D^{1 / 2}$-norm of $\mathscr{L}_{\omega}$ for the unsymmetric case are not discussed in the literature.

In this paper we will investigate these problems under the assumption that in (1.1) the matrix $A$ has the special form (1.7) with $D_{1}$ and $D_{2}$ symmetric and positive definite and $K^{T}=-H$. We will obtain some results similar to those for the symmetric case.

In the next section we review some properties for skew-symmetric matrices required for their application in the later sections. In Section 3 we construct the basis of eigenvectors of the associated Jacobi matrix $B$ which is similar to a skewsymmetric matrix. In Section 4 we will show that the Jordan canonical form of $\mathscr{L}_{\omega_{b}}$ is not a diagonal matrix, but has only $q$ principal vectors of grade 2 associated with $\omega_{b}-1$, the eigenvalues of $\mathscr{L}_{\omega_{b}}$. Here, $q$ is the multiplicity of the eigenvalue $i \bar{\mu}(=i \rho(B))$ of $B$. Hence, $\left\|\mathscr{L}_{\omega_{b}}^{m}\right\|$, the spectral norm, behaves like $m \cdot \rho\left(\mathscr{L}_{\omega_{b}}\right)^{m-1}$ rather than $\rho\left(\mathscr{L}_{\omega_{b}}\right)^{m}$.

In Section 5 we will determine the $D^{1 / 2}$-norm of $\mathscr{L}_{\omega}$ and point out that if $\bar{\mu}=$ $\rho(B) \geq 1$, then $\left\|\mathscr{L}_{\omega_{b}}^{m}\right\|_{D^{1 / 2}} \geq 1$. However, in Section 6 , we will show that for any $\bar{\mu}>0$, for $m$ large enough, $\left\|\mathscr{L}_{\omega_{b}}^{m}\right\|_{D^{1 / 2}}<1$. Eventually, $\left\|\mathscr{L}_{\omega_{b}}^{m}\right\|_{D^{1 / 2}}$ converges to zero, though considerably more slowly than $\rho\left(\mathscr{L}_{\omega_{b}}^{m}\right)$.

In this paper, almost all the notations used are the same as those adopted by Young [7], and all our work is based on the theory of Young [7].
2. Seme Properties of Skew-Symmetric Matrices. Let $A \in R^{n \times n}$ and

$$
\begin{equation*}
A^{T}=-A . \tag{2.1}
\end{equation*}
$$

It is well known that $A$ has the following properties:
(a) All diagonal elements of $A$ are zero.
(b) $A$ has either zero eigenvalues or purely imaginary eigenvalues, that is, any eigenvalue $\mu$ of $A$ has the form

$$
\begin{equation*}
\mu=i \xi \tag{2.2}
\end{equation*}
$$

Here $\xi$ is real. Also, $-\mu=-i \xi$ is an eigenvalue of $A$.
(c) $A$ is a normal matrix, that is,

$$
\begin{equation*}
A^{T} A=A A^{T} \tag{2.3}
\end{equation*}
$$

(d) $A$ is unitarily similar to a diagonal matrix.

All the above properties are easy to prove and can be found in any textbook of linear algebra, e.g., see [2].
3. The Eigenvectors of the Associated Jacobi Iteration Matrix B. Consider $A$ in (1.1) to have the special form

$$
A=\left(\begin{array}{cc}
D_{1} & -H  \tag{3.1}\\
-K & D_{2}
\end{array}\right)
$$

where $D_{1}\left(\in R^{r \times r}\right)$ and $D_{2}\left(\in R^{s \times s}\right)$ are symmetric positive definite and

$$
\begin{equation*}
H^{T}=-K \tag{3.2}
\end{equation*}
$$

the associated Jacobi iterative matrix $B$ has the form

$$
B=\left(\begin{array}{cc}
0 & F  \tag{3.3}\\
G & 0
\end{array}\right)
$$

Here,

$$
\begin{equation*}
G=D_{2}^{-1} K, \quad F=D_{1}^{-1} H \tag{3.4}
\end{equation*}
$$

Because $D_{1}$ and $D_{2}$ are positive definite, we can choose symmetric and positive definite matrices $D_{1}^{1 / 2}$ and $D_{2}^{1 / 2}$ such that

$$
\begin{equation*}
D_{1}^{1 / 2} D_{1}^{1 / 2}=D_{1}, \quad D_{2}^{1 / 2} D_{2}^{1 / 2}=D_{2} \tag{3.5}
\end{equation*}
$$

If we write

$$
D=\left(\begin{array}{cc}
D_{1}^{1 / 2} & 0  \tag{3.6}\\
0 & D_{2}^{1 / 2}
\end{array}\right)
$$

then we have

$$
D^{1 / 2} B D^{-1 / 2}=\left(\begin{array}{cc}
0 & D_{1}^{-1 / 2} H D_{2}^{-1 / 2}  \tag{3.7}\\
D_{2}^{-1 / 2} K D_{1}^{-1 / 2} & 0
\end{array}\right)
$$

Hence $B$ is similar to a skew-symmetric matrix, and thus unitarily similar to a diagonal matrix.

In this section we will construct a basis of eigenvectors for $B$. From (3.3) we have

$$
B^{2}=\left(\begin{array}{cc}
F G & 0  \tag{3.8}\\
0 & G F
\end{array}\right)
$$

Evidently, $B^{2}$ is also similar to a diagonal matrix and, in fact, the $(r \times r)$ matrix $F G$ and the ( $s \times s$ ) matrix $G F$ are also similar to diagonal matrices, where $r+s=n$, the order of the matrix $A$. Also note that $F G$ and $G F$ have nonpositive eigenvalues. Let the $p$ eigenvectors of $F G$ associated with the nonzero eigenvalues $\nu_{1}, \nu_{2}, \ldots, \nu_{p}$ be $\xi_{1}, \xi_{2}, \ldots, \xi_{p}$, i.e.,

$$
\begin{equation*}
F G \xi_{j}=\nu_{j} \xi_{j}, \quad j=1,2, \ldots, p \tag{3.9}
\end{equation*}
$$

If we let

$$
\begin{equation*}
\eta_{j}=G \xi_{j}, \quad j=1,2, \ldots, p, \tag{3.10}
\end{equation*}
$$

then $\eta_{j} \neq 0$, and $\eta_{j}$ is an eigenvector of $G F$ associated with $\nu_{j}$, i.e.,

$$
\begin{equation*}
G F \eta_{j}=\nu_{j} \eta_{j}, \quad j=1,2, \ldots, p \tag{3.11}
\end{equation*}
$$

Moreover, since the $\xi_{j}, j=1,2, \ldots, p$, are linearly independent, then so are the $\eta_{j}$, $j=1,2, \ldots, p$, since

$$
\sum_{j=1}^{p} c_{j} \eta_{j}=0
$$

implies that

$$
0=F\left(\sum_{j=1}^{p} c_{j} \eta_{j}\right)=\sum_{j=1}^{p} \nu_{j} c_{j} \xi_{j}=0
$$

and hence the $c_{j}, j=1,2, \ldots, p$, vanish because of the linear independence of the $\xi_{j}, j=1,2, \ldots, p$. Evidently, there can be no more than $p$ eigenvectors of $G F$ associated with nonzero eigenvalues; otherwise, there would be more than $p$ linearly independent eigenvectors of $F G$ associated with the nonzero eigenvalues. Thus, we have

$$
\begin{equation*}
p \leq \min \{r, s\} \tag{3.12}
\end{equation*}
$$

Since $\nu_{j}<0, j=1,2, \ldots, p$, if we let

$$
\begin{equation*}
\mu_{j}=i\left|\nu_{j}\right|^{1 / 2}, \quad x_{j}=\mu_{j} \xi_{j}, \quad y_{j}=\eta_{j}, \quad v_{j}=\binom{x_{j}}{y_{j}}, \quad j=1,2, \ldots, p \tag{3.13}
\end{equation*}
$$

where $i^{2}=-1$, then using (3.9), (3.11) and (3.13), we have

$$
\begin{equation*}
B v_{j}=\binom{F y_{j}}{G x_{j}}=\binom{\nu_{j} \xi_{j}}{\mu_{j} \eta_{j}}=\binom{\mu_{j} \mu_{j} \xi_{j}}{\mu_{j} \eta_{j}}=\mu_{j} v_{j}, \quad j=1,2, \ldots, p \tag{3.14}
\end{equation*}
$$

Notice that $\mu_{j}, j=1,2, \ldots, p$, have positive imaginary parts.
Let us now define for $j=p+1, p+2, \ldots, 2 p$

$$
\begin{equation*}
x_{j}=x_{j-p}, \quad y_{j}=-y_{j-p}, \quad v_{j}=\binom{x_{j}}{y_{j}}, \quad \mu_{j}=-\mu_{j-p} \tag{3.15}
\end{equation*}
$$

Evidently, we have

$$
\begin{equation*}
B v_{j}=\mu_{j} v_{j}, \quad j=p+1, p+2, \ldots, 2 p \tag{3.16}
\end{equation*}
$$

If we let $F G x=0$, where $x \neq 0$, then by (3.4) we have

$$
D_{1}^{-1} H D_{2}^{-1} K x=0, \quad \text { or } \quad H D_{2}^{-1} K x=0
$$

Thus, we have

$$
\begin{equation*}
-H D_{2}^{-1 / 2} D_{2}^{-1 / 2} H^{T} x=0 \quad \text { or } \quad\left(D_{2}^{-1 / 2} H^{T} x\right)^{*}\left(D_{2}^{-1 / 2} H^{T} x\right)=0 \tag{3.17}
\end{equation*}
$$

Here * stands for the conjugate transpose of a matrix. Hence from (3.17) we have $D_{2}^{-1 / 2} H^{T} x=0$, or $D_{2}^{-1} H^{T} x=G x=0$. Therefore, we have that if $F G x=0$, where $x \neq 0$, then

$$
\begin{equation*}
B\binom{x}{0}=\binom{0}{G x}=0 \tag{3.18}
\end{equation*}
$$

Thus, if the eigenvectors of $F G$ associated with the eigenvalue zero are $x_{2 p+1}$, $x_{2 p+2}, \ldots, x_{p+r}$, then the vectors

$$
\begin{equation*}
v_{j}=\binom{x_{j}}{0}, \quad j=2 p+1,2 p+2, \ldots, p+r \tag{3.19}
\end{equation*}
$$

are eigenvectors of $B$ associated with the eigenvalue zero. Similarly, if the eigenvectors of $G F$ associated with the eigenvalue zero are $y_{p+r+1}, y_{p+r+2}, \ldots, y_{s+r}$, then the vectors

$$
\begin{equation*}
v_{j}=\binom{0}{y_{j}}, \quad j=p+r+1, p+r+2, \ldots, n=r+s \tag{3.20}
\end{equation*}
$$

are eigenvectors of $B$ associated with the eigenvalue zero.
We have thus constructed a basis of eigenvectors for $B$,

$$
\begin{equation*}
v_{j}=\binom{x_{j}}{y_{j}}, \quad j=1,2, \ldots, n \tag{3.21}
\end{equation*}
$$

and moreover, we have by (3.14) and (3.16)

$$
\begin{equation*}
G x_{j}=\mu_{j} y_{j}, \quad F y_{j}=\mu_{j} x_{j}, \quad j=1,2, \ldots, n . \tag{3.22}
\end{equation*}
$$

We also have

$$
\begin{aligned}
& \text { for } \mu_{j}=i\left|\nu_{j}\right|^{1 / 2}, \quad v_{j}=\binom{x_{j}}{y_{j}}, \quad j=1,2, \ldots, p ; \\
& \text { for } \mu_{j}=-i\left|\nu_{j}\right|^{1 / 2}, \quad v_{j}=\binom{x_{j-p}}{-y_{j-p}}, \quad j=p+1, \ldots, 2 p ; \\
& \text { for } \mu_{j}=0, \quad v_{j}=\binom{x_{j}}{0}, \quad j=2 p+1, \ldots, r+p ; \\
& \text { for } \mu_{j}=0, \quad v_{j}=\binom{0}{y_{j}}, \quad j=p+r+1, \ldots, n .
\end{aligned}
$$

4. The Principal Vectors of $\mathscr{L}_{\omega}$. We now seek the eigenvectors and principal vectors of $\mathscr{L}_{\omega}$ for $\omega \neq 0$. Because $A$ has the form of (3.1), from (3.3) we have

$$
L=\left[\begin{array}{cc}
0 & 0  \tag{4.1}\\
G & 0
\end{array}\right], \quad U=\left[\begin{array}{cc}
0 & F \\
0 & 0
\end{array}\right]
$$

Thus we have

$$
\begin{align*}
\mathscr{L}_{\omega} & =(I-\omega L)^{-1}((1-\omega) I+\omega U) \\
& =\left[\begin{array}{cc}
I_{1} & 0 \\
-\omega G & I_{2}
\end{array}\right]^{-1}\left[\begin{array}{cc}
(1-\omega) I_{1} & \omega F \\
0 & (1-\omega) I_{2}
\end{array}\right] \\
& =\left[\begin{array}{cc}
I_{1} & 0 \\
\omega G & I_{2}
\end{array}\right]\left[\begin{array}{cc}
(1-\omega) I_{1} & \omega F \\
0 & (1-\omega) I_{2}
\end{array}\right]  \tag{4.2}\\
& =\left[\begin{array}{cc}
(1-\omega) I_{1} & \omega F \\
\omega(1-\omega) G & \omega^{2} G F+(1-\omega) I_{2}
\end{array}\right],
\end{align*}
$$

where $I_{1}$ and $I_{2}$ are identity matrices of the same sizes as $D_{1}$ and $D_{2}$, respectively. For each nonzero eigenvalue $\mu$ of $B$, let $\lambda_{+}^{1 / 2}$ and $\lambda_{-}^{1 / 2}$ be the roots of

$$
\begin{equation*}
\lambda+\omega-1=\omega \mu \lambda^{1 / 2} . \tag{4.3}
\end{equation*}
$$

Since

$$
B\binom{x}{y}=\mu\binom{x}{y},
$$

the vectors

$$
\begin{equation*}
w=\binom{x}{\lambda_{+}^{1 / 2} y}, \quad z=\binom{x}{\lambda_{-}^{1 / 2} y} \tag{4.4}
\end{equation*}
$$

are the eigenvectors of $\mathscr{L}_{\omega}$, since by (4.2), (3.22) and (4.3) we have

$$
\begin{align*}
\mathscr{L}_{\omega}\binom{x}{\lambda_{+}^{1 / 2} y} & =\binom{\left(1-\omega+\omega \mu \lambda_{+}^{1 / 2}\right) x}{\left[\omega \mu\left(1-\omega+\omega \mu \lambda_{+}^{1 / 2}\right)+(1-\omega) \lambda_{+}^{1 / 2}\right] y}  \tag{4.5}\\
& =\lambda_{+} w
\end{align*}
$$

and

$$
\begin{equation*}
\mathscr{L}_{\omega} z=\lambda_{-} z . \tag{4.6}
\end{equation*}
$$

If we let

$$
\begin{equation*}
v=\binom{x}{y}, \quad \hat{v}=\binom{x}{-y} \tag{4.7}
\end{equation*}
$$

then we have

$$
\begin{align*}
& w=\frac{1}{2}(v+\hat{v})+\frac{1}{2} \lambda_{+}^{1 / 2}(v-\hat{v}) \\
& z=\frac{1}{2}\left(1+\lambda_{+}^{1 / 2}\right) v+\frac{1}{2}\left(1-\lambda_{+}^{1 / 2}\right) \hat{v}  \tag{4.8}\\
&z+\hat{v})+\frac{1}{2} \lambda_{-}^{1 / 2}(v-\hat{v})=\frac{1}{2}\left(1+\lambda_{-}^{1 / 2}\right) v+\frac{1}{2}\left(1-\lambda_{-}^{1 / 2}\right) \hat{v}
\end{align*}
$$

If $\lambda_{-}^{1 / 2} \neq \lambda_{+}^{1 / 2}$, then $w$ and $z$ are linearly independent. But for $\omega \mu \neq 0$, the discriminant $\omega^{2} \mu^{2}-4(\omega-1)$ of (4.3) does not vanish unless

$$
\begin{equation*}
\omega^{2}|\mu|^{2}+4(\omega-1)=0 \tag{4.9}
\end{equation*}
$$

On the other hand, if (4.9) holds and if $\omega \mu \neq 0$, then $\lambda_{+}^{1 / 2}=\lambda_{-}^{1 / 2}=\lambda^{1 / 2}=\omega \mu / 2 \neq$ 0 , and $w$ and $z$ are not linearly independent. Notice that in this case,

$$
\begin{equation*}
\lambda_{+}=\lambda_{-}=\lambda=\omega^{2} \mu^{2} / 4=\omega-1 \tag{4.10}
\end{equation*}
$$

If we let

$$
\begin{equation*}
\hat{z}=\frac{1}{2} \cdot \frac{1}{\lambda^{1 / 2}}\binom{0}{y} \tag{4.11}
\end{equation*}
$$

then we have

$$
\begin{align*}
\mathscr{L}_{\omega} \hat{z} & =\left(\begin{array}{cc}
(1-\omega) I_{1} & \omega F \\
\omega(1-\omega) G & \omega^{2} G F+(1-\omega) I_{2}
\end{array}\right) \cdot \frac{1}{2 \lambda^{1 / 2}}\binom{0}{y} \\
& =\frac{1}{2 \lambda^{1 / 2}}\binom{\omega F y}{\omega^{2} G F y+(1-\omega) y}=\frac{1}{2 \lambda^{1 / 2}}\binom{\omega \mu x}{\omega \mu^{2} y+(1-\omega) y} \\
& =\frac{1}{2 \lambda^{1 / 2}}\binom{\omega \mu x}{\left[\omega^{2} \mu^{2}+1-\omega-\lambda+\lambda\right] y} \\
& =\frac{1}{2 \lambda^{1 / 2}}\binom{2 \cdot \frac{\omega \mu}{2} x}{\left[\omega^{2} \mu^{2}+1-\omega-(\omega-1)\right] y}+\lambda\binom{0}{\frac{1}{2 \lambda^{1 / 2}} y}  \tag{4.12}\\
& =\frac{1}{2 \lambda^{1 / 2}}\binom{2 \cdot \lambda^{1 / 2} x}{[4(\omega-1)+2(1-\omega)] y}+\lambda \hat{z}=\frac{1}{2 \lambda^{1 / 2}}\binom{2 \lambda^{1 / 2} x}{2(\omega-1) y}+\lambda \hat{z} \\
& =\binom{x}{\lambda^{1 / 2} y}+\lambda \hat{z}=w+\lambda \hat{z} .
\end{align*}
$$

Hence, $\hat{z}$ is a principal vector of grade 2. Moreover, we have

$$
\begin{equation*}
w=\frac{1}{2}\left(1+\lambda^{1 / 2}\right) v+\frac{1}{2}\left(1-\lambda^{1 / 2}\right) \hat{v}, \quad \hat{z}=\frac{1}{4} \cdot \frac{1}{\lambda^{1 / 2}} v-\frac{1}{4} \cdot \frac{1}{\lambda^{1 / 2}} \hat{v} . \tag{4.13}
\end{equation*}
$$

Thus $w$ and $\hat{z}$ are linearly independent.
If we let, for $\omega \mu \neq 0$,

$$
\begin{gathered}
w_{j}=\binom{x_{j}}{\left(\lambda_{j}\right)_{+}^{1 / 2} y_{j}}, \quad j=1,2, \ldots, p \\
w_{j+p}=\left\{\begin{array}{c}
\binom{x_{j}}{\left(\lambda_{j}\right)^{1 / 2} y_{j}}, \quad j=1,2, \ldots, p, \text { if } \omega^{2}\left|\mu_{j}\right|^{2}+4(\omega-1) \neq 0, \\
\frac{1}{2} \cdot \frac{1}{\left(\lambda_{j}\right)_{+}^{1 / 2}}\binom{0}{y_{j}}, \quad j=1,2, \ldots, p, \text { if } \cdot \omega^{2}\left|\mu_{j}\right|^{2}+4(\omega-1)=0 \\
w_{j}=v_{j}, \quad j=2 p+1,2 p+2, \ldots, n,
\end{array}\right.
\end{gathered}
$$

then we can easily prove that $w_{j}, j=1,2, \ldots, n$, are linearly independent and hence form a basis of the $n$-dimensional complex vector space $\mathbf{C}^{n}$. Therefore, the matrix whose columns are the $w_{j}$ reduces $\mathscr{L}_{\omega}$ to Jordan canonical form. Thus we have proved:

THEOREM 1. If the matrix $A$ has the form of (3.1) and $i \bar{\mu}=i \rho(B)$ is an eigenvalue of multiplicity $q$ of $B$ of (3.3), then the Jordan canonical form of $\mathscr{L}_{\omega_{b}}$ has $n-2 q(1 \times 1)$ sub-Jordan blocks and $q(2 \times 2)$ sub-Jordan blocks which correspond to the eigenvalue $\omega_{b}-1$.

Notice that if $q=1$, then the Jordan canonical form of $\mathscr{L}_{\omega_{b}}$ has one nondiagonal element.

From Theorem 3.1 of [5, p. 65] and Theorem 3-7.1 of [7, p. 85] we have

$$
\begin{equation*}
\left\|\mathscr{L}_{\omega_{b}}^{m}\right\| \sim J\left(\mathscr{L}_{\omega_{b}}\right) \cdot m \cdot \rho\left(\mathscr{L}_{\omega_{b}}\right)^{m-1} \tag{4.15}
\end{equation*}
$$

Here, $\|A\|$ is the spectral norm of the matrix $A$ and $J\left(\mathscr{L}_{\omega_{b}}\right)$ is the Jordan condition number of the matrix $\mathscr{L}_{\omega_{b}}$, defined by Young [7, p. 85] and given by

$$
\begin{equation*}
J\left(\mathscr{L}_{\omega_{b}}\right)=\inf _{V \in S_{1}} \kappa(V) \tag{4.16}
\end{equation*}
$$

where $\kappa(V)$ is the spectral condition number of the matrix $V$ and $S_{1}$ the set of all matrices such that

$$
\begin{equation*}
V^{-1} \mathscr{L}_{\omega_{b}} V=J \tag{4.17}
\end{equation*}
$$

Here, $J$ is the Jordan canonical form of the matrix $\mathscr{L}_{\omega_{b}}$.
5. Determination of $\left\|\mathscr{L}_{\omega}\right\|_{D^{1 / 2}}$. Let

$$
\begin{equation*}
\hat{\mathscr{L}}_{\omega}=D^{1 / 2} \mathscr{L}_{\omega} D^{-1 / 2} \tag{5.1}
\end{equation*}
$$

then from (4.2) and (3.6) we have

$$
\hat{\mathscr{L}}_{\omega}=\left(\begin{array}{cc}
(1-\omega) I_{1} & \omega D_{1}^{1 / 2} F D_{2}^{-1 / 2}  \tag{5.2}\\
\omega(1-\omega) D_{2}^{1 / 2} G D_{1}^{-1 / 2} & \omega^{2} D_{2}^{1 / 2} G F D_{2}^{-1 / 2}+(1-\omega) I_{2}
\end{array}\right) .
$$

If we let

$$
\begin{align*}
& \hat{F}=D_{1}^{1 / 2} F D_{2}^{-1 / 2}=D_{1}^{-1 / 2} H D_{2}^{-1 / 2}, \\
& \hat{G}=D_{2}^{1 / 2} G D_{1}^{-1 / 2}=D_{2}^{-1 / 2} K D_{1}^{-1 / 2}, \tag{5.3}
\end{align*}
$$

then

$$
\begin{equation*}
\hat{G}^{T}=-\hat{F} \tag{5.4}
\end{equation*}
$$

and

$$
\hat{\mathscr{L}}_{\omega}=\left(\begin{array}{cc}
(1-\omega) I_{1} & \omega \hat{F}  \tag{5.5}\\
\omega(1-\omega) \hat{G} & \omega^{2} \hat{G} \hat{F}+(1-\omega) I_{2}
\end{array}\right)
$$

Hence, $\hat{\mathscr{L}}_{\omega}$ is the SOR iterative matrix corresponding to the matrix

$$
\hat{A}=D^{1 / 2} A D^{-1 / 2}=\left(\begin{array}{cc}
I_{1} & -\hat{F}  \tag{5.6}\\
-\hat{G} & I_{2}
\end{array}\right)
$$

with the associated Jacobi iterative matrix

$$
\hat{B}=\left(\begin{array}{cc}
0 & \hat{F}  \tag{5.7}\\
\hat{G} & 0
\end{array}\right)=D^{1 / 2} B D^{-1 / 2}
$$

Therefore, $\rho(\hat{B})=\rho(B)$, and $\omega_{b}$ is the same for $\hat{A}$ as for $A$. Moreover, if we let $\mathscr{L}_{\omega}[A]$ stand for the SOR iterative matrix associated with the matrix $A$ and $D[A]$ for the diagonal block of the matrix $A$, then we have

$$
\begin{equation*}
\left\|\mathscr{L}_{\omega}^{m}[A]\right\|_{D[A]^{1 / 2}}=\left\|\hat{\mathscr{L}}_{\omega}^{m}[A]\right\|=\left\|\mathscr{L}_{\omega}^{m}[\hat{A}]\right\|=\left\|\mathscr{L}_{\omega}^{m}[\hat{A}]\right\|_{D[\hat{A}]^{1 / 2}} . \tag{5.8}
\end{equation*}
$$

Thus, it is sufficient to assume $A$ of (3.1) with $D_{1}=I_{1}$ and $D_{2}=I_{2}$. Otherwise, we consider $\hat{A}$ (5.6). Notice that when $D_{1}=I_{1}$ and $D_{2}=I_{2}$ then $F=H, G=K$, and $F^{T}=-G$.

Since

$$
\begin{equation*}
\left\|\mathscr{L}_{\omega}^{T}\right\|_{D^{1 / 2}}=\left\|\mathscr{L}_{\omega}\right\|_{D^{1 / 2}}=\left[\rho\left(\mathscr{L}_{\omega} \mathscr{L}_{\omega}^{*}\right)\right]^{1 / 2} \tag{5.9}
\end{equation*}
$$

according to the expression (4.2) for $\mathscr{L}_{\omega}$ we first study the eigenvalues of products of matrices of the form

$$
\left(\begin{array}{cc}
a_{11}(F G) & a_{12}(F G) F  \tag{5.10}\\
a_{21}(G F) G & a_{22}(G F)
\end{array}\right)
$$

where $a_{11}$ and $a_{12}$ are polynomials in $F G$ and $a_{21}$ and $a_{22}$ polynomials in $G F$. By an analogy with Theorem 7-2.1 of Young [7, p. 239] we have:

Theorem 2. If $B$ is a matrix of the form (3.3), then
(a) The matrix

$$
Q=\left(\begin{array}{cc}
a_{11}(F G) & a_{12}(F G) F  \tag{5.11}\\
a_{21}(G F) G & a_{22}(G F)
\end{array}\right)
$$

is nonsingular if

$$
\begin{equation*}
\tau\left(B^{2}\right)=a_{11}\left(B^{2}\right) a_{22}\left(B^{2}\right)-a_{21}\left(B^{2}\right) a_{12}\left(B^{2}\right) B^{2} \tag{5.12}
\end{equation*}
$$

in nonsingular. Moreover, $\tau\left(B^{2}\right)$ is nonsingular if and only if for each eigenvalue $\mu$ of $B$ the matrix

$$
R(\mu)=\left(\begin{array}{cc}
a_{11}\left(\mu^{2}\right) & a_{12}\left(\mu^{2}\right) \mu  \tag{5.13}\\
a_{21}\left(\mu^{2}\right) \mu & a_{22}\left(\mu^{2}\right)
\end{array}\right)
$$

is nonsingular.
(b) Let

$$
G_{m}=\prod_{k=m}^{1}\left(\begin{array}{cc}
a_{11}^{(k)}(F G) & a_{12}^{(k)}(F G) F  \tag{5.14}\\
a_{21}^{(k)}(G F) G & a_{22}^{k}(G F)
\end{array}\right)^{\nu_{k}}
$$

where for each $k, \nu_{k}= \pm 1$. It is assumed that for any $k$ the matrix

$$
\begin{equation*}
\tau^{(k)}\left(B^{2}\right)=a_{11}^{(k)}\left(B^{2}\right) a_{22}^{(k)}\left(B^{2}\right)-a_{21}^{(k)}\left(B^{2}\right) a_{12}^{(k)}\left(B^{2}\right) B^{2} \tag{5.15}
\end{equation*}
$$

is nonsingular for $\nu_{k}=-1$. For each eigenvalue $\mu$ of $B$, let

$$
M_{m}(\mu)=\prod_{k=m}^{1}\left(\begin{array}{cc}
a_{11}^{(k)}\left(\mu^{2}\right) & a_{12}^{(k)}\left(\mu^{2}\right) \mu  \tag{5.16}\\
a_{21}^{(k)}\left(\mu^{2}\right) \mu & a_{22}^{(k)}\left(\mu^{2}\right)
\end{array}\right)^{\nu_{k}}
$$

If $\mu$ is a nonzero eigenvalue of $B$ and if $\lambda$ is an eigenvalue of $M_{m}(\mu)$, then $\lambda$ is an eigenvalue of $G_{m}$. If $\mu=0$ is an eigenvalue of $B$, then at least one of the eigenvalues of $M_{m}(0)$ is an eigenvalue of $G_{m}$.
(c) If $\lambda$ is an eigenvalue of $G_{m}$, then there exists an eigenvalue $\mu$ of $B$ such that $\lambda$ is an eigenvalue of $M_{m}(\mu)$.

Notice that although the matrix $B$ considered here and the matrix $B$ considered in Theorem 7-2.1 of Young [7, p. 239] are not the same type of matrices-the former is similar to a skew-symmetric matrix and the latter a symmetric matrixthe statement of these two theorems are the same and the proofs are also the same. Hence the proof of Theorem 2 is omitted.

From (4.2) we have

$$
\begin{align*}
\mathscr{L}_{\omega}^{*} & =\mathscr{L}_{\omega}^{T}=\left(\begin{array}{cc}
(1-\omega) I_{1} & 0 \\
-\omega G & (1-\omega) I_{2}
\end{array}\right)\left(\begin{array}{cc}
I_{1} & -\omega F \\
0 & I_{2}
\end{array}\right) \\
& =\left(\begin{array}{cc}
(1-\omega) I_{1} & -\omega(1-\omega) F \\
-\omega G & \omega^{2} G F+(1-\omega) I_{2}
\end{array}\right) \tag{5.17}
\end{align*}
$$

Thus, from (4.2) and (5.7), $\mathscr{L}_{\omega}$ and $\mathscr{L}_{\omega}^{*}$ have the required form for the applicability of Theorem 2. From Theorem 2 we know that the eigenvalues of $\mathscr{L}_{\omega} \mathscr{L}_{\omega}^{*}$ are the same as the eigenvalues of $M(\omega, \mu) M^{*}(\omega, \mu)$, where

$$
\begin{align*}
M(\omega, \mu) & =\left(\begin{array}{cc}
1 & 0 \\
\omega \mu & 1
\end{array}\right)\left(\begin{array}{cc}
1-\omega & \omega \mu \\
0 & 1-\omega
\end{array}\right)  \tag{5.18}\\
& =\left(\begin{array}{cc}
1-\omega & \omega \mu \\
(1-\omega) \omega \mu & \omega^{2} \mu^{2}+1-\omega
\end{array}\right) .
\end{align*}
$$

If we notice $\bar{\mu}=-\mu$ (here $\mu$ is purely imaginary), and if we let

$$
M(\omega, \mu) M^{*}(\omega, \mu)=\left(\begin{array}{ll}
m_{11} & m_{12}  \tag{5.19}\\
m_{21} & m_{22}
\end{array}\right)
$$

then

$$
\begin{aligned}
m_{11} & =(1-\omega)^{2}-\omega^{2} \mu^{2} \\
m_{12} & =\omega^{3} \mu^{3}+\omega \mu(1-\omega)-(1-\omega)^{2} \cdot \omega \mu=\mu \omega^{2}\left[1-\omega+\omega \mu^{2}\right] \\
m_{21} & =-m_{12}=-\mu \omega^{2}\left[1-\omega+\omega \mu^{2}\right] \\
m_{22} & =\omega^{4} \mu^{4}+(1-\omega)^{2}+2 \omega^{2} \mu^{2}(1-\omega)-\omega^{2} \mu^{2}(1-\omega)^{2} \\
& =\omega^{4} \mu^{4}+(1-\omega)^{2}+\omega^{2} \mu^{2}(1-\omega)(1+\omega)
\end{aligned}
$$

Since

$$
\begin{align*}
m_{11}+m_{22} & =2(1-\omega)^{2}+\omega^{4} \mu^{4}-\omega^{4} \mu^{2} \\
m_{12} m_{21} & =-\mu^{2} \omega^{4}\left[(1-\omega)^{2}+2 \omega \mu^{2}(1-\omega)+\omega^{2} \mu^{4}\right] \\
m_{22} m_{11} & =-\omega^{6} \mu^{6}-2 \omega^{4} \mu^{4} \cdot \omega \cdot(1-\omega)-\omega^{2} \mu^{2}(1-\omega)^{2} \cdot \omega^{2}+(1-\omega)^{4}  \tag{5.20}\\
& =-\omega^{4} \mu^{2}\left[\omega^{2} \mu^{4}+2 \omega \mu^{2}(1-\omega)+(1-\omega)^{2}\right]+(1-\omega)^{4} \\
& =m_{21} m_{12}+(1-\omega)^{4}
\end{align*}
$$

we have

$$
\begin{equation*}
m_{22} m_{11}-m_{21} m_{12}=(1-\omega)^{4} \tag{5.21}
\end{equation*}
$$

Thus, if we let

$$
\begin{equation*}
\operatorname{det}\left(\lambda I-M(\omega, \mu) M^{*}(\omega, \mu)\right)=\lambda^{2}-T\left(\mu^{2}\right) \lambda+c=0 \tag{5.22}
\end{equation*}
$$

then

$$
\begin{equation*}
T\left(\mu^{2}\right)=2(1-\omega)^{2}+\omega^{4} \mu^{4}-\omega^{4} \mu^{2}, \quad c=(1-\omega)^{4} \tag{5.23}
\end{equation*}
$$

Notice that $\mu^{2} \leq 0$, so that $T\left(\mu^{2}\right)$ is an increasing function of $|\mu|$. Therefore, by Lemma 6-2.9 of Young [7, p. 186], it follows that for a given $\omega$, the largest value of the root radius of (5.22) is assumed for $\mu=i \rho(B)$ (or $\mu=-i \rho(B)$ ). From (5.22) and (5.23) we know that if $\lambda$ satisfies (5.22), then $t=\lambda^{1 / 2}$ satisfies

$$
\begin{equation*}
t^{2}-(\omega-1)^{2}=\omega^{2}|\mu|\left(1+|\mu|^{2}\right)^{1 / 2} t \tag{5.24}
\end{equation*}
$$

Note $|\mu|=\rho(B)$, and if we let $\bar{\mu}=\rho(B)$ and

$$
\begin{equation*}
d=\bar{\mu}\left(1+\bar{\mu}^{2}\right)^{1 / 2} \tag{5.25}
\end{equation*}
$$

then by Lemma 6-2.1 of Young [7, p. 171] the root radius of (5.24) is less than unity if and only if we have

$$
|\omega-1|<1 \quad \text { and } \quad \omega^{2} d<1-(\omega-1)^{2}=\omega(2-\omega)
$$

or, equivalently,

$$
\begin{equation*}
0<\omega<\min \left\{2, \frac{2}{1+d}\right\}=\frac{2}{1+d} \quad(\bar{\mu}>0) \tag{5.26}
\end{equation*}
$$

Thus we have proved
ThEOREM 3. If $A$ has the form (3.1) with $D_{1}$ and $D_{2}$ symmetric and positive definite and $H$ and $K$ satisfying (3.2), then $\left\|\mathscr{L}_{\omega}\right\|_{D^{1 / 2}}<1$ if and only if $\omega$ satisfies (5.26). Moreover, we have

$$
\begin{equation*}
\left\|\mathscr{L}_{\omega}\right\|_{D^{1 / 2}}=\frac{\omega^{2} d+\sqrt{\omega^{4} d^{2}+4(1-\omega)^{2}}}{2} \tag{5.27}
\end{equation*}
$$

We now determine the minimum value of $\left\|\mathscr{L}_{\omega}\right\|_{D^{1 / 2}}$. If we let

$$
f(\omega)=\omega^{2} d+\sqrt{\omega^{4} d^{2}+4(1-\omega)^{2}}
$$

then the derivative of $f(\omega)$ is given by

$$
f^{\prime}(\omega)=2 \omega d+\left[4 \omega^{3} d^{2}+8(\omega-1)\right] / 2 \sqrt{\omega^{4} d^{2}+4(1-\omega)^{2}}
$$

Assume $f^{\prime}(\omega)=0$; then

$$
\begin{equation*}
-\omega d \sqrt{\omega^{4} d^{2}+4(1-\omega)^{2}}=\omega^{3} d^{2}+2(\omega-1) \tag{5.28}
\end{equation*}
$$

Notice that (5.28) means

$$
\begin{equation*}
g(\omega)=\omega^{3} d^{2}+2(\omega-1)<0 \tag{5.29}
\end{equation*}
$$

By Descartes' rule we know that $g(\omega)$ has only one positive root $\omega_{u}$. Thus, if $\omega \in\left(0, \omega_{u}\right)$, then (5.29) holds. Moreover, if $\omega \geq \omega_{u}$, we have

$$
\begin{equation*}
f^{\prime}(\omega)>0 \tag{5.30}
\end{equation*}
$$

If $0<\omega<\omega_{u}$, and from (5.28), we have

$$
\begin{equation*}
\omega^{2} d^{2}+\omega-1=0 . \tag{5.31}
\end{equation*}
$$

Evidently, the positive root $\omega_{+}$of (5.31) is given by

$$
\begin{equation*}
\omega_{+}=\left[-1+\sqrt{1+4 d^{2}}\right] / 2 d^{2}=\frac{2}{1+\sqrt{1+4 d^{2}}} \tag{5.32}
\end{equation*}
$$

One can examine

$$
\begin{equation*}
\omega_{+}<\min \left\{2, \frac{2}{1+d}\right\}=\frac{2}{1+d} . \tag{5.33}
\end{equation*}
$$

Thus, we obtain

$$
f^{\prime}(\omega) \begin{cases}<0 & \text { if } 0<\omega<\omega_{+} \\ =0 & \text { if } \omega=\omega_{+} \\ >0 & \text { if } \omega>\omega_{+}\end{cases}
$$

because we have

$$
\begin{equation*}
\omega_{+}<\omega_{u} \tag{5.34}
\end{equation*}
$$

In fact, if $\omega_{+} \geq \omega_{u}$, then from (5.31) and (5.32) we have $\omega^{2} d^{2}+\omega-1<0$ for $0<\omega<\omega_{u}$. Thus we can prove $f^{\prime}(\omega)<0$ for $0<\omega<\omega_{u}$. Owing to the continuity property of $f^{\prime}(\omega)$, we have $f^{\prime}\left(\omega_{u}\right) \leq 0$, which contradicts (5.30). Hence (5.34) holds. We have now proved the following theorem.

ThEOREM 4. Under the assumptions of Theorem 3 we have

$$
\left\|\mathscr{L}_{\omega_{+}}\right\|_{D^{1 / 2}}<\left\|\mathscr{L}_{\omega}\right\|_{D^{1 / 2}} \quad \text { for } \omega \neq \omega_{+} .
$$

Here, $\omega_{+}$is given by (5.32).
It is important to note that from (1.9), (5.25), and (5.26) we have

$$
\begin{array}{ll}
\omega_{b}<\frac{2}{1+d} & \text { if } \bar{\mu}=\rho(B)<1 \\
\omega_{b} \geq \frac{2}{1+d} & \text { if } \bar{\mu} \geq 1
\end{array}
$$

Thus, when $\bar{\mu}<1$, we also have $\left\|\mathscr{L}_{\omega_{b}}\right\|_{D^{1 / 2}}<1$.
In fact we have proved
ThEOREM 5. Under the assumptions of Theorem 3 we have

$$
\left\|\mathscr{L}_{\omega_{b}}\right\|_{D^{1 / 2}}<1 \quad \text { if and only if } \bar{\mu}=\rho(B)<1 .
$$

But when $\bar{\mu} \geq 1$, we have $\left\|\mathscr{L}_{\omega_{b}}\right\|_{D^{1 / 2}} \geq 1$. However, in the next section, we will prove that for any $\bar{\mu}>0,\left\|\mathscr{L}_{\omega_{b}}^{m}\right\|_{D^{1 / 2}}<1$ if $m$ is large enough, and that $\operatorname{limit}_{m \rightarrow \infty}\left\|\mathscr{L}_{\omega_{b}}^{m}\right\|=0$.
6. Determination of $\left\|\mathscr{L}_{\omega_{b}}^{m}\right\|_{D^{1 / 2}}$. In this section we continue with the theory of Young [7, Chapter 7] to investigate $\left\|\mathscr{L}_{\omega_{b}}^{m}\right\|_{D^{1 / 2}}$. From the discussion of the last section it is sufficient to consider $A$ (3.1) with $D_{1}=I_{1}$ and $D_{2}=I_{2}$. Since the eigenvalues of $\mathscr{L}_{\omega_{b}}^{m}\left(\mathscr{L}_{\omega_{b}}^{m}\right)^{*}$ are the same as those of $M^{m}\left(\omega_{b}, \mu\right)\left[M^{m}\left(\omega_{b}, \mu\right)\right]^{*}$, where $M\left(\omega_{b}, \mu\right)$ is given by (5.18), we first develop an expression for $M^{m}(\omega, \mu)$. If we define the polynomials $S_{0}(\mu), S_{1}(\mu), \ldots$ by the recursion formula

$$
\begin{align*}
& S_{k}(\mu)=\omega \mu S_{k-1}(\mu)+(1-\omega) S_{k-2}(\mu), \quad k \geq 2  \tag{6.1}\\
& S_{0}(\mu)=1, \quad S_{1}(\mu)=\omega \mu
\end{align*}
$$

then by a result of Young [7, p. 248] we have

$$
M(\omega, \mu)^{m}=\left(\begin{array}{cc}
(1-\omega) S_{2 m-2} & S_{2 m-1}  \tag{6.2}\\
(1-\omega) S_{2 m-1} & S_{2 m}
\end{array}\right)
$$

Notice that $\mu$ is purely imaginary. By (6.1) one can see that $S_{2 k}(\mu)$ are real and $S_{2 k+1}$ purely imaginary. Also from the result of Young [7, p. 249, Eq. (4.7)] we have

$$
S_{k}(\mu)=\sum_{j=0}^{k} \alpha_{1}^{k-j} \alpha_{2}^{j}= \begin{cases}\frac{\alpha_{1}^{k+1}-\alpha_{2}^{k+1}}{\alpha_{1}-\alpha_{2}} & \text { if } \alpha_{1} \neq \alpha_{2}  \tag{6.3}\\ (k+1) \alpha^{k} & \text { if } \alpha_{1}=\alpha_{2}\end{cases}
$$

Here, $\alpha_{1}$ and $\alpha_{2}$ are the solution of the quadratic equation

$$
\begin{equation*}
\alpha^{2}-\omega \mu \alpha+\omega-1=0 \tag{6.4}
\end{equation*}
$$

Now we prove that if $\omega=\omega_{b}$ and $r=\left(1-\omega_{b}\right)$ then

$$
\begin{align*}
& S_{k}(i \bar{\mu})=S_{k}(i \rho(B))=(i)^{k}\left(r^{1 / 2}\right)^{k} \cdot(k+1),  \tag{6.5}\\
& \max _{\substack{\mu=i \beta \\
-\bar{\mu} \leq \beta \leq \bar{\mu}}}\left|S_{k}(\mu)\right|=\left|S_{k}(i \bar{\mu})\right|=(k+1)\left(r^{1 / 2}\right)^{k} . \tag{6.6}
\end{align*}
$$

Let $\mu=i \beta$; then the roots $\alpha_{1}$ and $\alpha_{2}$ of (6.4) are given by

$$
\alpha_{1,2}=\left[i \beta \omega_{b} \pm \sqrt{-\beta^{2} \omega_{b}^{2}-4\left(\omega_{b}-1\right)}\right] / 2=i\left[\beta \omega_{b} \pm \sqrt{\left.\beta^{2} \omega_{b}^{2}+4\left(\omega_{b}-1\right)\right]} / 2 .\right.
$$

By (1.9) we have

$$
\bar{\mu}^{2} \omega_{b}^{2}+4\left(\omega_{b}-1\right)=0 .
$$

Thus, if $\beta=\bar{\mu}$, then $\alpha_{1}=\alpha_{2}=i \bar{\mu} \omega_{b} / 2=i r^{1 / 2}$. Hence (6.5) follows from (6.3). If $|\beta| \leq \bar{\mu}$, then $\beta^{2} \omega_{b}^{2}+4\left(\omega_{b}-1\right) \leq 0$. Therefore, $\left|\alpha_{1}\right|=\left|\alpha_{2}\right|=\left(1-\omega_{b}\right)^{1 / 2}=r^{1 / 2}$. Again by (6.3), (6.6) follows.

From (6.2) we have

$$
\begin{align*}
M^{m} & (\omega, \mu)\left[M^{m}(\omega, \mu)\right]^{*} \\
& =\left(\begin{array}{cc}
(1-\omega) S_{2 m-2} & S_{2 m-1} \\
(1-\omega) S_{2 m-1} & S_{2 m}
\end{array}\right)\left(\begin{array}{cc}
(1-\omega) S_{2 m-2} & -(1-\omega) S_{2 m-1} \\
-S_{2 m-1} & S_{2 m}
\end{array}\right)  \tag{6.7}\\
& =\left(\begin{array}{cc}
(1-\omega)^{2} S_{2 m-2}^{2}-S_{2 m-1}^{2} & S_{2 m} S_{2 m-1}-(1-\omega)^{2} S_{2 m-1} S_{2 m-2} \\
(1-\omega)^{2} S_{2 m-1} S_{2 m-2}-S_{2 m} S_{2 m-1} & -(1-\omega)^{2} S_{2 m-1}^{2}+S_{2 m}^{2}
\end{array}\right) .
\end{align*}
$$

Evidently, the characteristic equation for $M^{m}\left(\omega_{b}, \mu\right)\left[M^{m}\left(\omega_{b}, \mu\right)\right]^{*}$ is

$$
\begin{equation*}
\lambda^{2}-T_{m}\left(\omega_{b}, \mu\right) \lambda+\Delta=0, \tag{6.8}
\end{equation*}
$$

where

$$
\begin{equation*}
T_{m}\left(\omega_{b}, \mu\right)=\left(1-\omega_{b}\right)^{2} S_{2 m-2}^{2}-S_{2 m-1}^{2}-\left(1-\omega_{b}\right)^{2} S_{2 m-1}^{2}+S_{2 m}^{2} \tag{6.9}
\end{equation*}
$$

and

$$
\begin{equation*}
\Delta=\operatorname{det}\left\{M^{m}\left(\omega_{b}, \mu\right)\left[M\left(\omega_{b}, \mu\right)\right]^{*}\right\}=r^{4 m}=\left(1-\omega_{b}\right)^{4 m} \tag{6.10}
\end{equation*}
$$

by (5.18). Since $\left[T_{m}\left(\omega_{b}, \mu\right)\right]^{2}-4 \Delta \geq 0$, because the eigenvalues of the Hermitian matrix $M^{m}\left(\omega_{b}, \mu\right)\left[M^{m}\left(\omega_{b}, \mu\right)\right]^{*}$ are real, it follows that for fixed $\Delta$ the modulus of the root of (6.8) is maximized when $\left|T_{m}\left(\omega_{b}, \mu\right)\right|$, considered as a function of $\mu$, is maximized. But, by (6.5) and (6.6), $\left|T_{m}\left(\omega_{b}, \mu\right)\right|$ is maximized when $\mu=i \bar{\mu}$, and we have

$$
\begin{align*}
\left|T_{m}\left(\omega_{b}, i \bar{\mu}\right)\right|= & r^{2} \cdot(2 m-1)^{2} r^{2 m-2}+r^{2 m-1} \cdot(2 m)^{2} \\
& +r^{2}(2 m)^{2} r^{2 m-1}+(2 m+1)^{2} r^{2 m}  \tag{6.11}\\
= & 2 r^{2 m}\left[1+2 m^{2}\left(\sqrt{r}+r^{-1 / 2}\right)^{2}\right] .
\end{align*}
$$

Thus, from (6.8), (6.10), and (6.11) we have

$$
\begin{equation*}
\left(\lambda-r^{2 m}\right)^{2}=4 m^{2}\left(r^{-1 / 2}+r^{1 / 2}\right)^{2} r^{2 m} \cdot \lambda \tag{6.12}
\end{equation*}
$$

and

$$
\begin{equation*}
\lambda-r^{2 m}=2 m\left(r^{-1 / 2}+r^{1 / 2}\right) r^{m} \lambda^{1 / 2} \tag{6.13}
\end{equation*}
$$

Hence we have proved the following
THEOREM 6. Under the assumptions of Theorem 3, we have

$$
\begin{align*}
\left\|\mathscr{L}_{\omega_{b}}^{m}\right\|_{D^{1 / 2}} & =r^{m}\left\{m\left(r^{-1 / 2}+r^{1 / 2}\right)+\left[m^{2}\left(r^{-1 / 2}+r^{1 / 2}\right)^{2}+1\right]^{1 / 2}\right\}  \tag{6.14}\\
& =F_{1}(m)
\end{align*}
$$

where

$$
\begin{equation*}
r=1-\omega_{b}, \quad \omega_{b}=\frac{2}{1+\sqrt{1+\bar{\mu}^{2}}}, \quad \bar{\mu}=\rho(B) \tag{6.15}
\end{equation*}
$$

From (6.14) we know that for any $\bar{\mu}=\rho(B)>0$

$$
\operatorname{limit}_{m \rightarrow \infty}\left\|\mathscr{L}_{\omega_{b}}^{m}\right\|_{D^{1 / 2}}=\operatorname{limit}_{m \rightarrow \infty} F_{1}(m)=0 .
$$

But, for values of $r$ close to unity, the function $F_{1}(m)$ increases initially before eventually decreasing. For $r$ close to unity we have

$$
\begin{align*}
\left\|\mathscr{L}_{\omega_{b}}^{m}\right\|_{D^{1 / 2}} & \sim 2 m r^{m}\left(r^{-1 / 2}+r^{1 / 2}\right) \\
& =2 m r^{m}\left(r^{-1} r^{1 / 2}+r^{-1} r \cdot r^{1 / 2}\right)  \tag{6.16}\\
& \sim 4 m r^{m-1}
\end{align*}
$$

On the other hand, we have

$$
\left\|\mathscr{L}_{\omega_{b}}^{m}\right\|=\left\|M^{m}\left(\omega_{b}, i \bar{\mu}\right)\right\| \sim m J\left(M\left(\omega_{b}, i \bar{\mu}\right)\right) r^{m-1}
$$

by Theorem 3-7.1 [7, p. 85]. Here, $J\left(M\left(\omega_{b}, i \bar{\mu}\right)\right)$ is the Jordan condition number of $M\left(\omega_{b}, i \bar{\mu}\right)$. But by [7, Theorem 3-8.1, p. 89] we have

$$
J\left(M\left(\omega_{b}, i \bar{\mu}\right)\right)=\omega_{b} \bar{\mu}+\left(1-\omega_{b}\right)_{\omega_{b}} \bar{\mu}=\omega_{b} \bar{\mu}\left(1+1-\omega_{b}\right)=2 r^{1 / 2}(1+r) \sim 4
$$

Hence,

$$
\begin{equation*}
\left\|\mathscr{L}_{\omega_{b}}^{m}\right\| \sim 4 m r^{m-1} \tag{6.17}
\end{equation*}
$$

Therefore, $\left\|\mathscr{L}_{\omega_{b}}^{m}\right\|$ behaves like $\left\|\mathscr{L}_{\omega_{b}}^{m}\right\|_{D^{1 / 2}}$.
Young [7, p. 255, Eq. (4.50)] has given $m_{0}$, the estimated number of iterations needed to reduce the $D^{1 / 2}$-norm of the error vector to a specified fraction $\varepsilon$ of the $D^{1 / 2}$-norm of the initial error vector as follows:

$$
\begin{align*}
m_{0} & =\log ((2 \nu / \varepsilon) \cdot \log (2 \nu / \varepsilon)) / \log (1 / r), \\
\nu & =\frac{r^{1 / 2}+r^{-1 / 2}}{\log (1 / r)} \tag{6.18}
\end{align*}
$$

Final Remarks. (a) Since $\left\|\mathscr{L}_{\omega_{b}}\right\|_{D^{1 / 2}} \geq 1$ if $\bar{\mu} \geq 1$, one should expect that it may be better to use $\omega=\omega_{+}$rather than $\omega=\omega_{b}$ in the initial steps. In this direction, an investigation is under way.
(b) By noting Theorem 6 and Theorem 7-4.1 of Young [7] one can find out that $\left\|\mathscr{L}_{\omega_{b}}^{m}\right\|_{D^{1 / 2}}$ for the nonsymmetric case and $\left\|\mathscr{L}_{\omega_{b}^{\prime}}^{m}\right\|_{D^{1 / 2}}$ for the symmetric and
positive definite case have the same expression in $m$ and $r$. The only diference is that for the former,

$$
\begin{equation*}
r=1-\omega_{b}=\frac{\rho^{2}(B)}{\left(1+\sqrt{1+\rho^{2}(B)}\right)}, \tag{6.19}
\end{equation*}
$$

and for the latter,

$$
\begin{equation*}
r=\omega_{b}^{\prime}-1=\frac{\rho^{2}(B)}{\left(1+\sqrt{1-\rho^{2}(B)}\right)} \tag{6.20}
\end{equation*}
$$

Especially for $m=1$, we have

$$
\begin{align*}
\left\|\mathscr{L}_{\omega_{b}}\right\|_{D^{1 / 2}} & =\left\|\mathscr{L}_{\omega_{b}}\right\|_{D^{1 / 2}}=r\left\{\left(r^{-1 / 2}+r^{1 / 2}\right)+\left[1+\left(r^{-1 / 2}+r^{1 / 2}\right)^{2}\right]^{1 / 2}\right\}  \tag{6.21}\\
& =r^{1 / 2}\left\{(1+r)+\left[r+(1+r)^{2}\right]^{1 / 2}\right\}=F(r)
\end{align*}
$$

It is clear that $F(r)$ is an increasing function of $r$. In fact one can prove
Lemma. Let $F(r)$ be given by (6.21). Then

$$
F(r)<1 \Leftrightarrow 0 \leq r<r_{0}=1 /(1+\sqrt{2})^{2} .
$$

By means of the above lemma we can give another proof of Theorem 5. In fact, it follows from (6.19), (6.21) and the above lemma that $\left\|_{\mathscr{L}_{w_{b}}}\right\|_{D^{1 / 2}}<1$ if and only if $\rho^{2}(B) /\left(1+\sqrt{1+\rho^{2}(B)}\right)<1 /(1+\sqrt{2})^{2}$, or equivalently, $\rho(B)<1$. Thus, Theorem 5 follows. However, we can give a similar result for the symmetric case. By noting (6.20), (6.21) and the above lemma, we have $\left\|\mathscr{L}_{\omega_{b}^{\prime}}\right\|_{D^{1 / 2}}<1$ if and only if $\rho^{2}(B) /\left(1+\sqrt{1-\rho^{2}(B)}\right)<1 /(1+\sqrt{2})^{2}$, or equivalently, $\rho(B)<1 / \sqrt{2}$. Thus we have proved

COROLLARY. If $A$ has the form (1.7) and is symmetric positive definite, then the $D^{1 / 2}$-norm of the corresponding optimum $S O R$ iterative matrix $\mathscr{L}_{\omega_{b}^{\prime}}$ is less than unity if and only if $\rho(B)<1 / \sqrt{2}$.

To our knowledge, the result of the above corollary is new. However, it should be noted that the result can be deduced from Theorem 7-3.1 of Young [7].

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