# An iteration method for solving the linear system $A x=b$ 

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#### Abstract

In this paper, based on a convergence splitting of the matrix $A$, we present an inner-outer iteration method for solving the linear system $A x=b$. We analyze the overall convergence of this method without any other restriction on its parameters. Moreover, we give the accelerated innerouter iteration method, and discuss how to apply the inner-outer iterations as a preconditioner for the krylov subspace methods. The inner-outer iteration method can also be used for the solution of $A X B=C$. Finally, several numerical examples are given to validate the performance of our proposed algorithms.


Key words: Inner-outer iterations; Convergence; Accelerate; Preconditioned

## 1. Introduction

Consider the solution of the following linear system

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

where $A \in R^{n \times n}$ is a nonsingular matrix, and $x, b \in R^{n}$.
The linear system (1.1) plays an important role in scientific and engineering computations. Many iteration methods have been proposed for solving (1.1) for decades, for example, the Jacobi, Gauss-Seidel and the Successive Overrelaxation (SOR) methods [1,6,10], etc. Krylov subspace methods, such as GMRES [7] and CG [5] have been considered as classical iteration methods for large and sparse linear systems. However, Krylov subspace methods may converge very slowly or not at all for many cases arising in certain computational mechanics and electronic device simulation problems. Hence preconditioning techniques [ $19,20,22,23$ ] have been widely used to improve the convergent behavior of Krylov subspace methods. For large non-Hermitian positive definite linear system, the HSS method $[2,4,8,11,15]$ has attracted more attention due to its promising performance and elegant mathematical properties.

In [3] the authors proposed an effective iteration method for PageRank computation [17,18,21] by using the inner-outer stationary iterations, which is very simple, and can be implemented and parallelized in a straightforward fashion, furthermore, this algorithm is an excellent combination of minimal memory requirements and fast convergence. In this paper, we will consider the

[^0]similar iteration method for the linear system (1.1). Moreover, the same idea can be used as a preconditioner for a nonstationary iteration method such as GMRES method.

The remainder of this paper is organized as follows: In section 2 we give the inner-outer iteration method and corresponding algorithm, discuss its convergence in detail. In section 3 we consider the overall convergence of the inner-outer iteration method and give a comparison result. In section 4 we present the accelerated inner-outer iteration method, and analyze how to use the inner-outer iteration as a preconditioning technique for nonstationary iteration methods. In section 5 we apply the inner-outer iteration method for solving the matrix equation $A X B=C$ and propose the corresponding algorithm. Several numerical examples are given in section 6 to illustrate the effectiveness of our proposed algorithms. Finally, we draw some conclusions in section 7.

## 2. The inner-outer iteration method for solving (1.1)

### 2.1 The inner-outer iteration method

Let $A=M-N$ be a convergent splitting, then we obtain the following linear system equivalent to (1.1):

$$
\begin{equation*}
(I-R) x=c \tag{2.1}
\end{equation*}
$$

where $R=M^{-1} N$ and $c=M^{-1} b$.
The inner-outer iteration method for solving (2.1) is expressed as follows:

$$
\begin{equation*}
(I-\alpha R) x_{k+1}=(1-\alpha) R x_{k}+c, k=0,1,2, \cdots \tag{2.2}
\end{equation*}
$$

with $0<\alpha<1$. Here we regard (2.2) as the outer iteration. It is difficult to obtain $x_{k+1}$ directly from (2.2) since it is equivalent to solving (1.1), then we approximate $x_{k+1}$ by the following inner iteration. Let $g=(1-\alpha) R x_{k}+c$, and we define the inner linear system as

$$
\begin{equation*}
(I-\alpha R) z=g \tag{2.3}
\end{equation*}
$$

then apply the inner iteration

$$
\begin{equation*}
z_{s+1}=\alpha R z_{s}+g, s=0,1,2, \cdots, l-1 \tag{2.4}
\end{equation*}
$$

where we take $z_{0}=x_{k}$ as the initial value and assign $z_{l}$ as the new $x_{k+1}$. If we have not converged to the desired solution of (2.1), then we repeat the above procedure, increasing $k$ and using $x_{k}$ as the new initial value $z_{0}$.

Assume the parameters $\eta$ and $\varepsilon$ are the tolerances of the inner and outer iterations, respectively, and we use the 2-norm of the residuals of the outer system (2.1) and the inner system (2.3) as stopping criterions to terminate the iterations. For the outer system (2.1) we have

$$
\left\|c-(I-R) x_{k+1}\right\|_{2}<\varepsilon
$$

and for the inner system (2.3) we apply

$$
\left\|g-(I-\alpha R) z_{s+1}\right\|_{2}<\eta .
$$

## Algorithm 1: The inner-outer iteration method

Input: $R, c, \alpha, \varepsilon, \eta$

Output: $x$
$x \leftarrow c$
$z \leftarrow R x$
while $\|c-x+z\|_{2} \geq \varepsilon$
$g \leftarrow(1-\alpha) z+c$
repeat
$x \leftarrow g+\alpha z$
$z \leftarrow R x$
until $\|g+\alpha z-x\|<\eta$
end while
0: $x \leftarrow z+c$
In this algorithm lines 1 and 2 initialize $x=c$ and $z=R x . g$ in (2.3) is obtained in line
4. The inner iteration defined in (2.4) is implemented from line 5 to line 8 . To terminate the Algorithm 1, line 3 tests the residual of the outer iteration (2.2), and for the inner iteration, the stopping criteria is in line 8. The vector $(I-\alpha R) z_{j+1}$ is given by $x-\alpha z$, since now $x$ holds $z_{j+1}$ and $z$ is computed by $R x$ in line 7 as well as in line 2 . Upon exit from the algorithm, $x$ is the desired approximation to the exact solution of the Eq.(2.1). In line 10 we execute a single iteration, since $z$ is already updated and we can use it as well.

Let $A=D-L-U$, where $D$ is the diagonal of $A,-L$ is the strictly lower triangular part of $A$, and $-U$ is the strictly upper triangular part of $A$, respectively. Then the matrix splitting of the AOR iteration method [24] for solving (1.1) is

$$
\begin{equation*}
M_{A}=\frac{1}{\omega}(D-\gamma L), \quad N_{A}=\frac{1}{\omega}((1-\omega) D+(\omega-\gamma) L+\omega U) \tag{2.5}
\end{equation*}
$$

where $\omega, \gamma$ are two real parameters with $\omega \neq 0$.
For different $\omega$ and $\gamma$, we can obtain the corresponding iteration methods:
(1) the Jacobi method: $\omega=1, \gamma=0$.
(2) the Gauss-Seidel method: $\omega=1, \gamma=1$.
(3) the SOR method: $\omega=\gamma$.

Then we can obtain the inner-outer iteration method based on the AOR splitting (2.5), and for different $\omega, \gamma$ we can get the corresponding inner-outer iteration methods.

In Example 1 the coefficient matrix $A$ of (1.1) is a nonsingular $M$ - matrix, then we need to discuss the convergence interval of the AOR iteration for this case, which will contribute to the choices of the parameters $\omega$ and $\gamma$. From [25] we can obtain the convergence theorem of the AOR iteration for a nonsingular $M-$ matrix $A$.
Theorem 2.1. Let $G_{A}=M_{A}^{-1} N_{A}$ and $J=D^{-1}(L+U)$ be the AOR iteration matrix and Jacobi iteration matrix for (1.1), respectively. If $A$ is a nonsingular $M$ - matrix, then

$$
\rho\left(G_{A}\right) \leq|1-\omega|+\omega \rho(J)<1
$$

for $0<\omega<\frac{2}{1+\rho(J)}$ and $0 \leq \gamma \leq \omega$.
Proof. This is a special case of Theorem 3.3 [25].
In order to improve the convergence behaviors of the iteration methods, many preconditioners [19,20,22,23] have been proposed for solving the linear system (1.1).

Firstly transforming (1.1) into the following preconditioned form

$$
\begin{equation*}
P A x=P b \tag{2.6}
\end{equation*}
$$

where $P \in R^{n \times n}$ is a preconditioner and nonsingular.
Let $P A=M_{p}-N_{p}$ be a convergence splitting and $M_{p}$ is nonsingular, then we have from (2.6)

$$
\begin{equation*}
\left(I-R_{p}\right) x=\bar{c} \tag{2.7}
\end{equation*}
$$

where $R_{p}=M_{p}^{-1} N_{p}, \bar{c}=M_{p}^{-1} \mathrm{~Pb}$. We can use the inner-outer iteration method (Algorithm 1) to solve (2.7), which is called the preconditioned inner-outer iteration method.
Remark 1. For the inner-outer iteration method we do not consider the case $\alpha=0$ or $\alpha=1$. If $\alpha=0$, the outer iteration (2.2) would simply lead back to the iteration sequence $x_{k+1}=R x_{k}+c$, and the inner iteration (2.4) is meaningless. For the case $\alpha=1$, the inner iteration (2.4) is also equivalent to the iteration sequence $x_{k+1}=R x_{k}+c$.

### 2.2 Convergence analysis of the inner-outer iteration method

In this section, we will discuss the convergence of the inner-outer iteration method. Here we will analyze the convergence of the iterations (2.2) and (2.4), respectively.
Theorem 2.2 [6]. The iteration sequence $x_{m+1}=R x_{m}+c$ converges to the solution of $A x=b$ for all starting vectors $x_{0}$ and for all $b$ if and only if $\rho(R)<1$.

The outer iteration (2.2) is associated with the matrix splitting

$$
I-R=M_{1}-N_{1}, M_{1}=I-\alpha R, N_{1}=(1-\alpha) R
$$

and the corresponding outer iteration matrix is

$$
\begin{equation*}
R_{1}=M_{1}^{-1} N_{1}=(1-\alpha)(I-\alpha R)^{-1} R . \tag{2.8}
\end{equation*}
$$

The inner iteration for solving (2.3) has the following matrix splitting

$$
M_{1}=I-\alpha R=M_{2}-N_{2}, M_{2}=I, \quad N_{2}=\alpha R,
$$

and the inner iteration matrix is

$$
\begin{equation*}
R_{2}=\alpha R \tag{2.9}
\end{equation*}
$$

Theorem 2.3. Given $\rho(R)<1$ and $0<\alpha<1$, then the outer iteration (2.2) is convergent, where $\rho(A)$ denotes the spectral radius of matrix $A$.
Proof. Let $\lambda_{i}$ be an eigenvalue of $R$, then $\left|\lambda_{i}\right| \leq \rho(R)<1$ from the assumption. Assume that $\mu_{i}$ is an eigenvalue of $R_{1}$ in (2.8), then

$$
\begin{equation*}
\mu_{i}=\frac{(1-\alpha) \lambda_{i}}{1-\alpha \lambda_{i}} \tag{2.10}
\end{equation*}
$$

Since $\left|\lambda_{i}\right| \leq \rho(R)<1$ and $0<\alpha<1$, then from (2.10) we have

$$
\begin{equation*}
\left|\mu_{i}\right|=\left|\frac{(1-\alpha) \lambda_{i}}{1-\alpha \lambda_{i}}\right| \leq \frac{(1-\alpha)|\rho(R)|}{1-\alpha|\rho(R)|}=\frac{1-\alpha}{|\rho(R)|^{-1}-\alpha}<1 . \tag{2.11}
\end{equation*}
$$

From Theorem 2.2 and (2.11) the outer iteration (2.2) is convergent.
Theorem 2.4. If $\rho(R)<1$ and $0<\alpha<1$, then the inner iteration (2.4) is convergent.
Proof. Let $\phi_{i}$ be an eigenvalue of $R_{2}$ in (2.9), then

$$
\begin{equation*}
\phi_{i}=\alpha \lambda_{i}, \tag{2.12}
\end{equation*}
$$

where $\lambda_{i}$ is an eigenvalue of $R$. Then from (2.12) we get

$$
\begin{equation*}
\left|\phi_{i}\right|=\left|\alpha \lambda_{i}\right| \leq \alpha \rho(R)<1 . \tag{2.13}
\end{equation*}
$$

From Theorem 2.2 and (2.13) we complete the proof.
Remark 2. Let $A=\bar{M}_{i}-\bar{N}_{i}(i=1,2)$ be two convergence splittings. If $\rho\left(\bar{M}_{1}^{-1} \bar{N}_{1}\right)<$ $\rho\left(\bar{M}_{2}^{-1} \bar{N}_{2}\right)<1$, then from Theorems 2.3 and 2.4 it is clear that the eigenvalues of the iteration matrices $\bar{R}_{1}$ and $\bar{R}_{2}$ based on $\left(\bar{M}_{1}, \bar{N}_{1}\right)$ have a smaller upper bound, then the inner-outer iteration method generated by $\left(\bar{M}_{1}, \bar{N}_{1}\right)$ maybe converges faster than that derived from $\left(\bar{M}_{2}, \bar{N}_{2}\right)$.

## 3. The overall convergence of the inner-outer iteration method

In this section, we will give the proof of the overall convergence for the inner-outer iteration method without considering the parameters $\varepsilon$ and $\eta$, which shows that the inner-outer iteration method converges linearly to the exact solution $x^{*}$ of the linear system (2.1).
Lemma 3.1 [6]. For all operator norms $\rho(R) \leq\|R\|$. For all $R$ and for all $\varepsilon>0$ there is an operator norm $\|R\|_{\star} \leq \rho(R)+\varepsilon$. The norm $\|\cdot\|_{\star}$ depends on both $R$ and $\varepsilon$, where $\rho(C)$ denotes the spectral radius of the matrix $C$.
Lemma 3.2 [6]. Let $\|\cdot\|$ satisfy $\|A B\| \leq\|A\| \cdot\|B\|$. Then $\|X\|<1$ implies that $I-X$ is invertible, $(I-X)^{-1}=\sum_{i=0}^{\infty} X^{i}$, and $\left\|(I-X)^{-1}\right\| \leq \frac{1}{1-\|X\|}$.

Firstly we rewrite the inner-outer iteration method as a two-stage iteration framework [17,32]:

$$
\left\{\begin{array}{l}
x_{k, 0}=x_{k}, x_{0}=c, x_{k+1}=x_{k, m_{k}}  \tag{3.1}\\
x_{k, j+1}=\alpha R x_{k, j}+(1-\alpha) R x_{k}+c, k=0,1,2, \cdots \\
j=0,1,2, \cdots, m_{k}-1
\end{array}\right.
$$

Theorem 3.1. Let $A=M-N$ be a convergence splitting and $0<\alpha<1$, and $m_{k}$ be the number of the inner iteration steps at the $k$-th outer iteration. Then the iteration sequence $\left\{x_{k}\right\}_{k=0}^{\infty}$ generated by (3.1) converges to the exact solution $x^{*}$ of (2.1), and faster than the iteration sequence derived from (2.1) for the same initial value $x_{0}$.
Proof. From (3.1) it follows that

$$
x_{k, j+1}=\left[(\alpha R)^{j+1}+(1-\alpha) \sum_{s=0}^{j}(\alpha R)^{s} R\right] x_{k}+\sum_{s=0}^{j}(\alpha R)^{s} c
$$

Then we have

$$
\begin{equation*}
x_{k+1}=H_{k} x_{k}+E_{k} c, \quad k=0,1,2, \cdots, \tag{3.2}
\end{equation*}
$$

where $H_{k}=(\alpha R)^{m_{k}}+(1-\alpha) \sum_{s=0}^{m_{k}-1}(\alpha R)^{s} R$ and $E_{k}=\sum_{s=0}^{m_{k}-1}(\alpha R)^{s}$.
Since $x^{*}$ is the exact solution of (1.1), then from (3.2) we obtain

$$
\begin{equation*}
x^{*}=H_{k} x^{*}+E_{k} c, \quad k=0,1,2, \cdots . \tag{3.3}
\end{equation*}
$$

Subtracting (3.3) from (3.2), then

$$
x_{k+1}-x^{*}=H_{k}\left(x_{k}-x^{*}\right)=\cdots=H_{k} H_{k-1} \cdots H_{0}\left(x_{0}-x^{*}\right), \quad k=0,1,2, \cdots
$$

and

$$
\begin{align*}
H_{k} & =(\alpha R)^{m_{k}}+(1-\alpha) \sum_{s=0}^{m_{k}-1}(\alpha R)^{s} R \\
& =(\alpha R)^{m_{k}}-\sum_{s=0}^{m_{k}-1}(\alpha R)^{s}[(I-R)-(I-\alpha R)] \\
& =(\alpha R)^{m_{k}}+\left(I-(\alpha R)^{m_{k}}\right)-\sum_{s=0}^{m_{k}-1}(\alpha R)^{s}(I-R)  \tag{3.4}\\
& =I-\sum_{s=0}^{m_{k}-1}(\alpha R)^{s}(I-R) .
\end{align*}
$$

Let $\lambda_{i}$ be an eigenvalue of $R$, then from(3.4) we have

$$
\begin{equation*}
\phi_{i}^{(k)}=1-\frac{\left(1-\lambda_{i}\right)\left(1-\left(\alpha \lambda_{i}\right)^{m_{k}}\right)}{1-\alpha \lambda_{i}} \tag{3.5}
\end{equation*}
$$

is an eigenvalue of $H_{k}$.
Since $0<\alpha<1$ and $\left|\lambda_{i}\right|<1$, then from (3.5) it is clear that

$$
\begin{align*}
\left|\phi_{i}^{(k)}\right| & =\left|1-\frac{\left(1-\lambda_{i}\right)\left(1-\left(\alpha \lambda_{i}\right)^{m_{k}}\right)}{1-\alpha \lambda_{i}}\right| \\
& =\left|\frac{\lambda_{i}\left(1-\alpha+\alpha^{m_{k}} \lambda_{i}^{m_{k}-1}-\left(\alpha \lambda_{i}\right)^{m_{k}}\right)}{1-\alpha \lambda_{i}}\right| \\
& <\frac{|1-\alpha|+\left|\alpha \lambda_{i}\right|^{m_{k}-1}\left|1-\lambda_{i}\right| \alpha}{\left|1-\alpha \lambda_{i}\right|}  \tag{3.6}\\
& <\frac{|1-\alpha|+\left|\alpha \lambda_{i}\right|^{m_{k}-1}(1+\rho(R)) \alpha}{\left|1-\alpha \lambda_{i}\right|} \\
& =\frac{|1-\alpha|}{\left|1-\alpha \lambda_{i}\right|}<1
\end{align*}
$$

as $m_{k} \rightarrow \infty$. Then $\rho\left(H_{k}\right)<1$.
Let $\psi=\max _{k}\left\{\rho\left(H_{k}\right)\right\}<1(k=0,1,2, \cdots)$ and $\nu_{i}$ be an eigenvalue of $H_{k} H_{k-1} \cdots H_{0}$, then we get

$$
\nu_{i}=\Pi_{s=0}^{k} \phi_{i}^{(s)}=\Pi_{s=0}^{k}\left(1-\frac{\left(1-\lambda_{i}\right)\left(1-\left(\alpha \lambda_{i}\right)^{m_{s}}\right)}{1-\alpha \lambda_{i}}\right),
$$

so

$$
\rho\left(H_{k} H_{k-1} \cdots H_{0}\right)=\max _{i}\left\{\left|\nu_{i}\right|\right\}=\max _{i}\left\{\Pi_{s=0}^{k}\left|\phi_{i}^{(s)}\right|\right\} \leq \rho\left(H_{k}\right) \rho\left(H_{k-1}\right) \cdots \rho\left(H_{0}\right) \leq \psi^{k+1}<1
$$

From Lemma 3.1 there exists an operator norm $\|\cdot\|_{\varrho}$ such that

$$
\left\|H_{k} H_{k-1} \cdots H_{0}\right\|_{\varrho} \leq \psi^{k+1}+\varepsilon
$$

for $\forall \varepsilon>0$. Then we have

$$
\begin{equation*}
\left\|x_{k+1}-x^{*}\right\|_{\varrho} \leq\left\|H_{k} H_{k-1} \cdots H_{0}\right\|_{\varrho}\left\|x_{0}-x^{*}\right\|_{\varrho} \leq\left(\psi^{k+1}+\varepsilon\right)\left\|x_{0}-x^{*}\right\|_{\varrho} . \tag{3.7}
\end{equation*}
$$

So from (3.7) we learn that the iteration sequence (3.1) converges to the exact solution $x^{*}$ as $\varepsilon \rightarrow 0$.

From (3.6) it is clear that

$$
\begin{aligned}
\left|\left(1-\alpha \lambda_{i}\right) \phi_{i}^{(k)}\right| & =\left|\lambda_{i}\left(1-\alpha+\alpha^{m_{k}} \lambda_{i}^{m_{k}-1}-\left(\alpha \lambda_{i}\right)^{m_{k}}\right)\right| \\
& =\left|\lambda_{i}\right| \cdot\left|1-\alpha+\alpha\left(\alpha \lambda_{i}\right)^{m_{k}-1}\left(1-\lambda_{i}\right)\right| \\
& <\left|\lambda_{i}\right|\left(|1-\alpha|+2 \alpha\left|\alpha \lambda_{i}\right|^{m_{k}-1}\right) \\
& =\left|\lambda_{i}\right| \cdot|1-\alpha| \\
& <\left|\lambda_{i}\right| \cdot\left|1-\alpha \lambda_{i}\right|
\end{aligned}
$$

with $m_{k} \rightarrow \infty$, then $\left|\phi_{i}^{(k)}\right|<\left|\lambda_{i}\right|$, so $\rho\left(H_{k}\right)<\rho(R)$ for $k=0,1,2, \cdots$.
Let $\left\{y_{k}\right\}_{k=0}^{\infty}$ be generated by the iteration sequence $y_{k}=R y_{k-1}+c$ from (2.1) with the same initial value $x_{0}$, and $x_{0}-x^{*}$ be an eigenvector of $R$ with eigenvalue $\lambda_{i}$ for $\left|\lambda_{i}\right|=\rho(R)$, then

$$
\begin{aligned}
\left\|x_{k+1}-x^{*}\right\| & =\left\|H_{k} H_{k-1} \cdots H_{0}\left(x_{0}-x^{*}\right)\right\| \\
& =\left\|\phi_{i}^{(k)} \phi_{i}^{(k-1)} \cdots \phi_{i}^{(0)}\left(x_{0}-x^{*}\right)\right\| \\
& \leq\left\|\rho\left(H_{k}\right) \rho\left(H_{k-1}\right) \cdots \rho\left(H_{0}\right)\left(x_{0}-x^{*}\right)\right\| \\
& <\left\|\rho(R) \rho(R) \cdots \rho(R)\left(x_{0}-x^{*}\right)\right\| \\
& =\left\|R^{k+1}\left(x_{0}-x^{*}\right)\right\|=\left\|y_{k+1}-x^{*}\right\|,
\end{aligned}
$$

we complete the proof.
Remark 3. From Theorem 3.1 we know that the conclusions are also held for an appropriate $m_{k}$. For example, let $\alpha=0.3$ and $\rho(R)=0.8$, then $\left|\alpha \lambda_{i}\right| \leq 0.24$, if $m_{k}=16$, then $\left|\alpha \lambda_{i}\right|^{m_{k}-1}=$ $\left|\alpha \lambda_{i}\right|^{15} \leq 5.04 \times 10^{-10}$.

## 4. The accelerated inner-outer iteration method and precondition

### 4.1 The accelerated inner-outer iteration method

In this section, we combine the Algorithm 1 with other solvers as an acceleration scheme, and obtain the accelerated inner-outer iteration method. At first we can use the inner-outer iterations, then switch to other solvers when the inner iterations converge fast.
Algorithm 2: The accelerated inner-outer iteration method
Input: $R, c, \alpha, \varepsilon, m$
Output: $x$
$x \leftarrow c$
$z \leftarrow R x$
while $\|c-x+z\|_{2} \geq \varepsilon$
$g \leftarrow(1-\alpha) z+c$
for $i=1: \kappa$
$x \leftarrow g+\alpha z$
$z \leftarrow R x$
end
If $i \leq \kappa, x=z+c, z=R x$; return
end while
. $x \leftarrow z+c$
In Algorithm 2 we use a for loop to count the number of the inner iterations. In line 9, we carry out the iteration sequences derived from the splitting $A=M-N$ with the latest $x$ as an initial value.

### 4.2 Preconditioning for nonstationary iteration methods

Let $A=M-N$ be a convergence splitting, then we can use $(I-\alpha R)^{-1} M^{-1}$ as a preconditioner for the linear system (1.1). Since the computation of the preconditioner is as difficult to solve (1.1), then we use the Neumann series approximation of $(I-\alpha R)^{-1} M^{-1}$ as a preconditioner. If $0<\alpha<1$ and $\rho(R)<1$, then from Lemmas 3.1 and 3.2 we have

$$
\begin{equation*}
(I-\alpha R)^{-1}=\sum_{i=0}^{\infty}(\alpha R)^{i} \tag{4.1}
\end{equation*}
$$

We can obtain an approximation of $(I-\alpha R)^{-1} M^{-1}$ by adopting $m$ terms from (4.1):

$$
\begin{equation*}
(I-\alpha R)^{-1} M^{-1} \approx\left(I+\alpha R+(\alpha R)^{2}+\cdots+(\alpha R)^{m}\right) M^{-1} \tag{4.2}
\end{equation*}
$$

Let $\lambda_{i}(i=1, \cdots n)$ be an eigenvalue of $R$ and $A=M-N$. If we precondition (1.1) with (4.2), then we have

$$
\begin{aligned}
Q & =\left(I+\alpha R+(\alpha R)^{2}+\cdots+(\alpha R)^{m}\right) M^{-1} A \\
& =\left(I+\alpha R+(\alpha R)^{2}+\cdots+(\alpha R)^{m}\right) M^{-1}(M-N) \\
& =\left(I+\alpha R+(\alpha R)^{2}+\cdots+(\alpha R)^{m}\right)(I-R)
\end{aligned}
$$

and

$$
\begin{equation*}
\tilde{\lambda}_{i}=\left(1-\lambda_{i}\right)\left(1+\alpha \lambda_{i}+\cdots+\left(\alpha \lambda_{i}\right)^{m}\right) \tag{4.3}
\end{equation*}
$$

where $\tilde{\lambda}_{i}(i=1, \cdots n)$ is an eigenvalue of the matrix $Q$.
In order to show the effectiveness of the preconditioner (4.2) on the krylov subspace methods, now we will analyze the distribution of the eigenvalues of $Q$.

From (4.3) we have

$$
\begin{align*}
\left|\tilde{\lambda}_{i}-1\right| & =\left|\frac{\left(1-\lambda_{i}\right)\left(1-\left(\alpha \lambda_{i}\right)^{m+1}\right)}{1-\alpha \lambda_{i}}-1\right| \\
& =\left|\frac{\left(\alpha \lambda_{i}\right)^{m+1}\left(\lambda_{i}-1\right)+\lambda_{i}(\alpha-1)}{1-\alpha \lambda_{i}}\right|  \tag{4.4}\\
& <\frac{2(\alpha \rho(R))^{m+1}+\rho(R)(1-\alpha)}{1-\alpha \rho(R)},
\end{align*}
$$

then all the eigenvalues of $Q$ are scattered in a circle centered at (1,0) with the radius $\frac{2(\alpha \rho(R))^{m+1}+\rho(R)(1-\alpha)}{1-\alpha \rho(R)}$. Let $\alpha=\rho(R)$, then from (4.4) we get

$$
\left|\tilde{\lambda}_{i}-1\right|<\frac{2 \rho(R)^{2 m+2}+\rho(R)(1-\rho(R))}{1-\rho(R)^{2}} .
$$

If $\rho(R)=0.7$ and $m=2$, then $\frac{2 \rho(R)^{2 m+2}+\rho(R)(1-\rho(R))}{1-\rho(R)^{2}}=0.8731$.
Remark 4. Let $A=\tilde{M}_{i}-\tilde{N}_{i}(i=1,2)$ be two convergence splittings and $\rho\left(\tilde{M}_{1}^{-1} \tilde{N}_{1}\right)<$ $\rho\left(\tilde{M}_{2}^{-1} \tilde{N}_{2}\right)$, then from (4.4) we obtain

$$
\frac{2\left(\alpha \rho\left(\tilde{R}_{1}\right)\right)^{m+1}+\rho\left(\tilde{R}_{1}\right)(1-\alpha)}{1-\alpha \rho\left(\tilde{R}_{1}\right)}<\frac{2\left(\alpha \rho\left(\tilde{R}_{2}\right)\right)^{m+1}+\rho\left(\tilde{R}_{2}\right)(1-\alpha)}{1-\alpha \rho\left(\tilde{R}_{2}\right)}
$$

so the preconditioned matrix $Q$ based on $\left(\tilde{M}_{1}, \tilde{N}_{1}\right)$ may be more effective.

## 5. The inner-outer iteration method for the matrix equation $A X B=C$

Firstly we will review some well-known definitions and lemmas, which can be found in $[1,6]$.
Let $X$ be an $m \times n$ matrix, then $\operatorname{vec}(X)$ is defined to be a column vector of size $m \cdot n$ made of the columns of $X$ stacked atop one another from left to right. Let $A$ and $B$ be $m \times n$ and $p \times q$ matrices, respectively. Then the Kronecker product $A \otimes B$ is the $(m \cdot p) \times(n \cdot q)$ matrix

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

Lemma 5.1. Let $A, B$ and $X$ be $m \times m, n \times n$, and $m \times n$ matrices, respectively. Then the following properties hold:

1. $\operatorname{vec}(A X)=\left(I_{n} \otimes A\right) \cdot \operatorname{vec}(X) ;$
2. $\operatorname{vec}(X B)=\left(B^{T} \otimes I_{m}\right) \cdot \operatorname{vec}(X)$.

Lemma 5.2. The following relations about Kronecker products hold:

1. Assume that the products $A \cdot C$ and $B \cdot D$ are well defined. Then $(A \otimes B) \cdot(C \otimes D)=$ $(A \cdot C) \otimes(B \cdot D) ;$
2. if $A$ and $B$ are invertible, then $(A \otimes B)^{-1}=A^{-1} \otimes B^{-1}$;
3. $(A \otimes B)^{T}=A^{T} \otimes B^{T}$.

Lemma 5.3. Let $\lambda(A)$ and $\mu(B)$ be the spectrums of $A \in R^{n \times n}$ and $B \in R^{m \times m}$, respectively, then

$$
\lambda(A \otimes B)=\left\{\lambda_{i} \mu_{j}: \lambda_{i} \in \lambda(A), \mu_{j} \in \mu(B), i=1,2, \cdots, n ; j=1,2, \cdots, m\right\} .
$$

In [12] the authors proposed a class of iteration methods for solving the matrix equation

$$
\begin{equation*}
A X B=C \tag{5.1}
\end{equation*}
$$

Since the matrix equation (5.1) is very important in matrix theory and applications, hence how to effectively solve this kind of equation has been under research recently, many iteration methods [9,13,14, 16] have been proposed for solving (5.1) and corresponding matrix equations.

Based on the matrix splittings of the matrix $A$ ro $B$, we can obtain the corresponding innerouter iteration method for solving (5.1). In this section we only consider the inner-outer iteration method for the matrix splittings of $B$.

Let $B^{T}=\hat{M}-\hat{N}$ be a convergent matrix splitting and $A$ be nonsingular, then we can get the equivalent form of (5.1) as follows:

$$
\begin{equation*}
X(I-\hat{R})=\hat{C} \tag{5.2}
\end{equation*}
$$

where $\hat{R}=\hat{N}^{T}\left(\hat{M}^{T}\right)^{-1}, \hat{C}=A^{-1} C\left(\hat{M}^{T}\right)^{-1}$.
The outer iteration for (5.2):

$$
\begin{equation*}
X_{k+1}(I-\beta \hat{R})=(1-\beta) X_{k} \hat{R}+\hat{C}, k=0,1,2, \cdots \tag{5.3}
\end{equation*}
$$

with $0<\beta<1$.
The inner iteration for (5.3):

$$
\begin{equation*}
Y_{t+1}=\beta Y_{t} \hat{R}+\hat{G}, t=0,1,2, \cdots, p-1 \tag{5.4}
\end{equation*}
$$

with $\hat{G}=(1-\beta) X_{k} \hat{R}+\hat{C}$. Here we take $Y_{0}=X_{k}$ as the initial value and assign $Y_{p}$ as the new $X_{k+1}$.

Similar to Algorithm 1, the inner-outer iteration method for (5.1) can be described as follows:
Algorithm 3: The inner-outer iteration method for (5.1)
Input: $\hat{R}, \hat{C}, \beta, \delta, \zeta$
Output: $X$
1: $X \leftarrow \hat{C}$
2: $Y \leftarrow X \hat{R}$
: while $\|\hat{C}-X+Y\|_{2} \geq \delta$
4: $\quad \hat{G} \leftarrow(1-\beta) Y+\hat{C}$
5: repeat
6: $\quad X \leftarrow \hat{G}+\beta Y$
7: $\quad Y \leftarrow X \hat{R}$
8: $\quad$ until $\|\hat{G}+\beta Y-X\|_{2}<\zeta$
9: end while
10: $X \leftarrow Y+\hat{C}$
Now we will analyze the convergence of the inner-outer iterations (5.3) and (5.4). By Lemma 5.1 we firstly transform the matrix equation (5.1) into the following form:

$$
\begin{equation*}
\left(B^{T} \otimes A\right) x=\bar{c} \tag{5.5}
\end{equation*}
$$

where $x$ and $\bar{c}$ are the forms of $\operatorname{vec}(X)$ and $\operatorname{vec}(C)$. From the matrix splitting $B^{T}=\hat{M}-\hat{N}$ and Lemma 5.2, we have

$$
\begin{equation*}
(I-\bar{R}) x=\hat{c} \tag{5.6}
\end{equation*}
$$

where $\bar{R}=\hat{M}^{-1} \hat{N} \otimes I$ and $\hat{c}=(\hat{M} \otimes A)^{-1} \bar{c}$.
The inner-outer iterations for (5.6):

$$
\begin{gather*}
(I-\beta \bar{R}) x_{k+1}=(1-\beta) \bar{R} x_{k}+\hat{c}, k=0,1,2, \cdots  \tag{5.7}\\
y_{s+1}=\beta \bar{R} y_{s}+\hat{g}, s=0,1,2, \cdots, p-1 \tag{5.8}
\end{gather*}
$$

with $\hat{g}=(1-\beta) \bar{R} x_{k}+\hat{c}$.
From Theorems 2.3, 2.4 and Lemma 5.3, we can get the following convergent theorems for (5.7) and (5.8) as well as (5.3) and (5.4), respectively.

Theorem 5.1. Given $\rho\left(\hat{M}^{-1} \hat{N}\right)<1$ and $0<\beta<1$, then the outer iteration (5.7) is convergent. Proof. Since $\rho\left(\hat{M}^{-1} \hat{N}\right)<1$, then from Lemma 5.3 we have $\rho(\bar{R})=\rho\left(\hat{M}^{-1} \hat{N} \otimes I\right)<1$. The proof is completed similar to that of Theorem 2.3.
Theorem 5.2. If $\rho\left(\hat{M}^{-1} \hat{N}\right)<1$ and $0<\beta<1$, then the inner iteration (5.8) is convergent.
Proof. From Lemma 5.3 and the assumptions we obtain $\rho(\bar{R})=\rho\left(\hat{M}^{-1} \hat{N} \otimes I\right)<1$, then we complete the proof according to Theorem 2.4.

Remark 5. If we rewrite the inner-outer iterations (5.7)-(5.8) as a two-stage iteration framework just like (3.1), we can also get the corresponding conclusions similar to Theorem 3.1.

## 6. Numerical results

In this section, we perform some numerical examples to illustrate the effectiveness of our algorithms in this paper. The numerical experiments are performed in Matlab R2010 on an Intel dual core processor ( $2.30 \mathrm{GHz}, 4 \mathrm{~GB}$ RAM). We use four iteration parameters to test these iteration methods, which are iteration step (denoted as IT), computing time in seconds (denoted as CPU), number of matrix-vectors (denoted as MV), and relative residual (denoted as RES) defined by $\frac{\left\|r_{k}\right\|_{2}}{\|b\|_{2}}$ with $r_{k}=b-A x_{k}$. Three sparse matrices $P$ that we will use in the following examples are listed in Table 1, where "Average Nonzeros" means the average number of the nonzero elements per row. For each matrix we will deal with it by some techniques, then each matrix $P$ is a column-stochastic matrix. Let $A=I-\varphi P$, where $I$ is an identity matrix and $0<\varphi<1$, then $A$ is a nonsingular $M$-matrix.

Table 1: Three test matrices for (1.1).

| Name | Size | Nonzeros | Average Nonzeros |
| :---: | :---: | :---: | :---: |
| Minnesota | $2,642 \times 2,642$ | 6,606 | 2.50 |
| Wb-cs-stanford | $9,914 \times 9,914$ | 36,854 | 3.71 |
| Stanford-Berkeley | $683,446 \times 683,446$ | $7,583,376$ | 11.0 |

Example 1. In this example, we use the AOR method, the inner-outer (IO) method (3.1) and the accelerated inner-outer (AIO) method (Algorithm 2) based on the AOR splitting (2.5) for solving (1.1), respectively, where $m_{k}=2$ and $\kappa=2$. The test matrix $P$ is Minnesota matrix (available from http://www.cise.ufl.edu/research/sparse/matrices/Gleich/index.html) and $\varphi=$ $0.95, b(i)=1(i=1, \cdots, n)$. All algorithms are terminated once the residual norms are below $10^{-8}$.

Numerical results are reported in Tables 2, 3. From Table 2 we find that AIO method perform better than the AOR and IO methods for different $\omega, \gamma$ in both iteration number and CPU time, but it needs more MV compared with other two methods except $\omega=1.8$ and $\gamma=1.7$. Moreover, all the iteration methods have more efficiency when $\omega, \gamma$ are near 1.5. From Table 3 we also get the similar conclusions for the SOR, IO and AIO methods, where the IO and AIO methods are derived from the matrix splitting (2.5) with $\omega=\gamma$.
Example 2. This example is devoted to comparing the iteration number, CPU time and MV of the IO method for different $\eta$ and $\alpha$, and the IO method is based on the Richardson iteration. The test matrix $P$ is the Stanford-Berkeley matrix of order 683,446 (available from http://cise.ufl.edu/research/sparse/matrices/SNAP/email-Enron.html) and $\varphi=0.95, b=(1-$

Table 2: Numerical results of the AOR, IO and AIO methods

| $\omega$ | $\gamma$ | $\rho(R)$ | Iteration method | IT | MV | CPU | RES |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| 1.8 | 1.7 | 0.8446 | The AOR method | 114 | 114 | 4.7011 | $8.79 \times 10^{-9}$ |
|  |  |  | The IO method | 32 | 64 | 1.5123 | $8.38 \times 10^{-9}$ |
|  |  |  | The AIO method | 23 | 69 | 1.2242 | $4.44 \times 10^{-9}$ |
| 1.6 | 1.5 | 0.7010 | The AOR method | 56 | 56 | 2.4827 | $7.29 \times 10^{-9}$ |
|  |  |  | The IO method | 30 | 60 | 1.5081 | $7.56 \times 10^{-9}$ |
|  |  |  | The AIO method | 23 | 69 | 1.2694 | $7.14 \times 10^{-9}$ |
| 1.5 | 1.4 | 0.7506 | The AOR method | 70 | 70 | 3.0767 | $8.59 \times 10^{-9}$ |
|  |  |  | The IO method | 40 | 80 | 1.8571 | $6.40 \times 10^{-9}$ |
|  |  |  | The AIO method | 30 | 90 | 1.5610 | $7.64 \times 10^{-9}$ |
| 1.2 | 1.1 | 0.9398 | The AOR method | 136 | 136 | 5.5520 | $9.19 \times 10^{-9}$ |
|  |  |  | The IO method | 76 | 152 | 3.5596 | $9.26 \times 10^{-9}$ |
|  |  |  | The AIO method | 57 | 171 | 2.9287 | $8.60 \times 10^{-9}$ |
| 0.9 | 0.8 | 0.9268 | The AOR method | 244 | 244 | 10.256 | $9.61 \times 10^{-9}$ |
|  |  |  | The IO method | 136 | 272 | 6.2308 | $9.64 \times 10^{-9}$ |
|  |  |  | The AIO method | 101 | 303 | 5.1983 | $9.67 \times 10^{-9}$ |

Table 3: Numerical results of the SOR, IO and AIO methods

| $\omega$ | $\gamma$ | $\rho(R)$ | Iteration method | IT | MV | CPU | RES |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| 1.7 | 1.7 | 0.7457 | The SOR method | 67 | 67 | 2.8872 | $8.22 \times 10^{-9}$ |
|  |  |  | The IO method | 32 | 64 | 1.5238 | $6.55 \times 10^{-9}$ |
|  |  |  | The AIO method | 20 | 60 | 1.0836 | $9.46 \times 10^{-9}$ |
| 1.5 | 1.5 | 0.6814 | The SOR method | 57 | 57 | 2.5109 | $8.66 \times 10^{-9}$ |
|  |  |  | The IO method | 32 | 64 | 1.5286 | $9.58 \times 10^{-9}$ |
|  |  |  | The AIO method | 24 | 72 | 1.3163 | $5.98 \times 10^{-9}$ |
| 1.2 | 1.2 | 0.8545 | The SOR method | 121 | 121 | 5.2262 | $9.02 \times 10^{-9}$ |
|  |  |  | The IO method | 68 | 136 | 3.1752 | $8.30 \times 10^{-9}$ |
|  |  |  | The AIO method | 51 | 153 | 2.7113 | $7.96 \times 10^{-9}$ |
| 0.9 | 0.9 | 0.9204 | The SOR method | 224 | 224 | 9.8298 | $9.80 \times 10^{-9}$ |
|  |  |  | The IO method | 125 | 250 | 5.8433 | $9.67 \times 10^{-9}$ |
|  |  |  | The AIO method | 90 | 270 | 4.8832 | $9.74 \times 10^{-9}$ |
| 0.8 | 0.8 | 0.9350 | The SOR method | 276 | 276 | 11.818 | $9.49 \times 10^{-9}$ |
|  |  |  | The IO method | 154 | 308 | 7.0680 | $9.26 \times 10^{-9}$ |
|  |  |  | The AIO method | 114 | 342 | 5.8866 | $9.72 \times 10^{-9}$ |

$\varphi) v / n$, where $v(i)=1(i=1, \cdots, n)$. All iterations are terminated once the residual norms are below $10^{-8}$.

From the discussion of $\eta$ in [3,21], a very strict $\eta$ may take a long time to compute the inner iteration in Algorithm 1, and just to implement a single outer iteration once a time. Setting $\eta$ very loose, however, may make the inner iteration do not sufficiently approximate the exact solution of (2.4), and this may slow the convergence rate overall.

From Table 4 we find that the IO method may converge faster for a smaller $\eta$, but take more computing time and MV for a given stopping criterion, for example $\eta=10^{-6}$. On the contrary, for a looser $\eta$ the IO method has a slower convergence than other four values of $\eta$, such as $\eta=10^{-2}$ or $\eta=10^{-3}$.

In table 5 we perform the IO method with different $\alpha$ and $\eta=10^{-4}$. From Table 5 we observe that the IO method is not sensitive to the choice of $\alpha$ in CPU time. When $\alpha$ is close to 1 , it needs less iteration number, but has more MV at the same time.

Table 4: Numerical results of the IO method for the Stanford-Berkeley matrix with different $\eta$.

| $\alpha=0.5$ | Iter | MV | CPU | RES |
| :--- | :---: | :---: | :---: | :---: |
| $\eta=10^{-2}$ | 422 | 422 | 20.382 | $9.11 \times 10^{-9}$ |
| $\eta=10^{-3}$ | 422 | 422 | 20.338 | $9.22 \times 10^{-9}$ |
| $\eta=10^{-4}$ | 410 | 440 | 20.685 | $9.34 \times 10^{-9}$ |
| $\eta=10^{-5}$ | 389 | 516 | 22.353 | $9.33 \times 10^{-9}$ |
| $\eta=10^{-6}$ | 367 | 663 | 25.668 | $9.14 \times 10^{-9}$ |

Table 5: Numerical results of the IO method for the Stanford-Berkeley matrix with different $\alpha$.

| $\eta=10^{-4}$ | Iter | MV | CPU | RES |
| :---: | :---: | :---: | :---: | :---: |
| $\alpha=0.2$ | 420 | 431 | 20.636 | $9.04 \times 10^{-9}$ |
| $\alpha=0.4$ | 414 | 438 | 20.711 | $9.11 \times 10^{-9}$ |
| $\alpha=0.5$ | 410 | 440 | 20.740 | $9.34 \times 10^{-9}$ |
| $\alpha=0.7$ | 401 | 445 | 20.489 | $9.13 \times 10^{-9}$ |
| $\alpha=0.9$ | 389 | 445 | 20.506 | $9.14 \times 10^{-9}$ |

Example 3: This example mainly illustrates the effectiveness of the preconditioner (4.2) for
the linear system (1.1) by using the GMRES method, where we choose $m=2$ and $\varphi=0.99$. The test matrices $P$ are the Minnesota matrix and Wb -cs-stanford matrix which has the order of 9,914 (available from http://www.cise.ufl.edu/research/sparse/matrices/Gleich/index.html.), and $b(i)=1(i=1, \cdots, n)$. In this example we use three preconditioners, which are the preconditioner (4.2) based on the Jacobi splitting and the Gauss-Seidel (GS) splitting, and the ILU preconditioner with no fill-in (ILU(0)), respectively.

Fig. 1 shows the distribution of the eigenvalues of the matrix $Q$ for the preconditioners based on the Jacobi and GS splittings for different $\alpha$, respectively. we find that the matrix $Q$ for the preconditioner(GS) has more eigenvalues concentrated around 1, which is more obvious for the larger $\alpha$, so the PGMRES method with the preconditioner(GS) is more efficient than the PGMRES method with the preconditioner(Jacobi), which is consistent with Remark 4. From Table 6 we learn that the PGMRES method with the three preconditioners all outperforms the GMRES method in both IT and CPU time, and the PGMRES(ILU(0)) performs best. From Table 7 we notice that the PGMRES(GS) has the same iteration number as the PGMRES(ILU(0)) for $\alpha=0.8$, while the PGMRES(ILU(0)) takes more CPU time than the PGMRES(GS) and PGMRES(Jacobi) for difficult $\alpha$, respectively.


Figure 1: The distribution of eigenvalues for the minnesota matrix

Example 4. In this example $A$ and $B$ are strictly diagonally dominant matrix, $A(i, i)=$ $6, A(i, i+1)=3, A(i+1, i)=-1, A(i, i+2)=-1$ and $B(i, i)=8, B(i, i+1)=-1, B(i+1, i)=$

Table 6: Numerical results for the minnesota matrix

| $\alpha$ | m | Iteration method | IT | CPU | RES |
| :---: | :---: | :---: | :---: | :---: | :---: |
| 0.2 | 2 | PGMRES(Jacobi) | 108 | 1.3747 | $8.55 \times 10^{-10}$ |
|  |  | PGMRES(GS) | 89 | 1.0794 | $9.83 \times 10^{-10}$ |
| 0.4 | 2 | PGMRES(Jacobi) | 91 | 1.1233 | $9.35 \times 10^{-10}$ |
|  |  | PGMRES(GS) | 74 | 0.8820 | $8.89 \times 10^{-10}$ |
| 0.6 | 2 | PGMRES(Jacobi) | 81 | 0.9832 | $9.09 \times 10^{-10}$ |
|  |  | PGMRES(GS) | 62 | 0.7342 | $7.91 \times 10^{-10}$ |
| 0.8 | 2 | PGMRES(Jacobi) | 76 | 0.9078 | $8.95 \times 10^{-10}$ |
|  |  | PGMRES(GS) | 53 | 0.6140 | $7.17 \times 10^{-10}$ |
|  |  | PGMRES(ILU(0)) | 42 | 0.4689 | $9.53 \times 10^{-10}$ |
|  |  | GMRES | 131 | 1.8129 | $8.76 \times 10^{-10}$ |

Table 7: Numerical results for the Wb-cs-stanford matrix

| $\alpha$ | m | Iteration method | IT | CPU | RES |
| :---: | :---: | :---: | :---: | :---: | :---: |
| 0.2 | 2 | PGMRES(Jacobi) | 121 | 1.7350 | $9.20 \times 10^{-10}$ |
|  |  | PGMRES(GS) | 79 | 1.5441 | $8.76 \times 10^{-10}$ |
| 0.4 | 2 | PGMRES(Jacobi) | 96 | 1.1428 | $8.99 \times 10^{-10}$ |
|  |  | PGMRES(GS) | 64 | 1.0750 | $9.08 \times 10^{-10}$ |
| 0.6 | 2 | PGMRES(Jacobi) | 79 | 0.8412 | $9.62 \times 10^{-10}$ |
|  |  | PGMRES(GS) | 53 | 0.7468 | $7.06 \times 10^{-10}$ |
| 0.8 | 2 | PGMRES(Jacobi) | 71 | 0.6610 | $9.36 \times 10^{-10}$ |
|  |  | PGMRES(GS) | 45 | 0.6469 | $6.34 \times 10^{-10}$ |
|  |  | PGMRES(ILU(0)) | 45 | 1.7448 | $7.76 \times 10^{-10}$ |
|  | GMRES | 159 | 2.0648 | $9.63 \times 10^{-10}$ |  |

$-1, B(i, i+2)=-1$ for $i=1,2, \cdots, n ; C_{i j}=1$ for $i, j=1,2, \cdots, n$.
We compare the IO method based on the Jacobi splitting with the Jacobi-type method and the HSS method, respectively, where we apply the IO method similar to (3.1) with $\beta=0.5$ and $m_{k}=2$. From Table 8 we find that the IO method is more efficient than the Jacobi-style method and HSS method in term of iteration number and CPU time, respectively.

Table 8: Numerical results for $A X B=C$

|  | Jacobi-type method |  |  | HSS method |  |  | The IO method $\left(m_{k}=2\right)$ |  |  |
| ---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| n | IT | CPU | RES | IT | CPU | RES | IT | CPU | RES |
| 200 | 97 | 0.1017 | $5.42 \times 10^{-10}$ | 62 | 0.6162 | $6.38 \times 10^{-10}$ | 29 | 0.0511 | $5.74 \times 10^{-10}$ |
| 500 | 100 | 1.6448 | $5.71 \times 10^{-10}$ | 65 | 9.5407 | $9.24 \times 10^{-10}$ | 30 | 1.0152 | $4.30 \times 10^{-10}$ |
| 800 | 100 | 6.1598 | $2.63 \times 10^{-10}$ | 67 | 16.021 | $7.08 \times 10^{-10}$ | 30 | 3.7903 | $5.44 \times 10^{-10}$ |
| 1000 | 101 | 11.277 | $2.94 \times 10^{-10}$ | 67 | 36.703 | $4.20 \times 10^{-10}$ | 30 | 6.5638 | $2.39 \times 10^{-10}$ |

Example 5. In this example $A$ is an $M$-matrix defined by

$$
A=\left[\begin{array}{cccccc}
10 & -1 & -2 & -1 & -3 & -1 \\
-3 & 11 & -2 & -1 & -3 & -2 \\
-1 & -2 & 15 & -5 & -2 & -3 \\
-3 & -4 & -1 & 14 & -2 & -2 \\
-3 & -5 & -1 & -3 & 16 & -1 \\
-1 & -2 & -3 & -4 & -2 & 16
\end{array}\right]
$$

and $b=\left[\begin{array}{llllll}1 & 1 & 1 & 1 & 1 & 1\end{array}\right]^{T}$.
We apply the preconditioned IO method for solving (1.1), which is based on the Jacobi and Gauss-Seidel splittings, respectively. We use the preconditioner (1.4) [23], and let $\alpha=$ $0.8, m_{k}=2$. From Tables 9,10 we find that the preconditioned IO method outperforms other three iteration methods in term of IT and CPU time. Moreover, the preconditioned IO method (GS) has more effectiveness than the preconditioned IO method (Jacobi), and performs best among all the iteration methods in this example.

## 7. Conclusion

In this paper, we give an inner-outer iteration method for solving the linear system (1.1). Furthermore, the inner-outer iterations can also be used as a preconditioner for the Krylov

Table 9: Numerical results of Example 5

| Iteration method | IT | CPU | RES |
| :--- | :---: | :---: | :---: |
| Jacobi method | 136 | 0.000819 | $8.77 \times 10^{-10}$ |
| Preconditioned Jacobi method | 113 | 0.000682 | $8.25 \times 10^{-10}$ |
| The IO method (Jacobi) | 82 | 0.000563 | $9.88 \times 10^{-10}$ |
| The preconditioned IO method (Jacobi) | 68 | 0.000336 | $9.04 \times 10^{-10}$ |

Table 10: Numerical results of Example 5

| Iteration method | IT | CPU | RES |
| :--- | :---: | :---: | :---: |
| Gauss-Seidel method | 68 | 0.000337 | $7.45 \times 10^{-10}$ |
| Preconditioned Gauss-Seidel method | 59 | 0.000247 | $7.92 \times 10^{-10}$ |
| The IO method (GS) | 39 | 0.000225 | $8.12 \times 10^{-10}$ |
| The preconditioned IO method (GS) | 32 | 0.000173 | $7.06 \times 10^{-10}$ |

subspace method such as GMRES method. We also apply the inner-outer iteration method for solving the matrix equation $A X B=C$, which is more efficient than the iteration methods [12]. Future work may further improve the IO iteration method and construct more effective algorithms for solving (1.1) and (5.1), respectively.

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