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## Shifted linear systems in electromagnetics.

Part II: Systems with multiple right-hand sides.

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#### Abstract

We consider the solution of multiply shifted linear systems for multiple right-hand sides. The coefficient matrix is symmetric, complex, and indefinite. The matrix is shifted by different multiples of the identity. Such problems arise in a number of applications, including the electromagnetic simulation in the development of microwave and mm-wave circuits and modules. The properties of microwave circuits can be described in terms of their scattering matrix which is extracted from the orthogonal decomposition of the electric field. We discretize the Maxwell's equations with orthogonal grids using the Finite Integration Technique (FIT). Some Krylov subspace methods have been used to solve systems with multiple right-hand sides. We use both the block-QMR method and a symmetric band Lanczos process based on coupled recurrences with polynomial preconditioning. We present a method for providing initial guesses to a linear solver both for systems with multiple shifts and for systems with multiple right-hand sides each with a different shift.


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## 1 Introduction

Today, electromagnetic simulation forms an indispensable tool in the development of microwave circuits. The description of the boundary of the computational domain has always been a key issue in bringing up efficiency of electromagnetic simulation. The Perfectly Matched Layer (PML) concept provides an excellent solution to this issue. However, the benefits of PML do not come for free. In the frequency-domain case, the material tensors worsen the numerical properties of the system of equations to be solved, which results in increased CPU time [20].
The subject under investigation are three-dimensional structures of arbitrary geometry which are connected to the remaining circuit by transmission lines. Ports are defined at the outer terminations of the transmission lines (see Fig. 1).


Figure 1: The basic structure under investigation

Calculating the excitations at the ports, one obtains eigenvalue problems and then large-scale systems of linear algebraic equations. In general, the computation of the eigenvalue problem and of the system of linear algebraic equations have to be done for several frequencies. Moreover, these linear equation problems have to be solved repeatedly for different right-hand sides. The number of right-hand sides depends on the number of ports and modes.

## 2 Scattering Matrix

The scattering matrix describes the structure in terms of the wave modes on the transmission line sections at the ports. We consider all exciting modes with amplitudes $a_{l}$ towards the discontinuity and all amplitudes $b_{l}$ outwards from the discontinuity (see Fig. 1). As example for the waves at the left port of Fig. 1 the transverse mode field at a cross-sectional plane $z$ is given by

$$
\begin{equation*}
\vec{E}_{t}(z)=\sum_{l=1}^{m^{(p)}} a_{l} \vec{E}_{t, l} e^{-\jmath k_{z_{l}} z}+\sum_{l=1}^{m^{(p)}} b_{l} \vec{E}_{t, l} e^{+\jmath k_{z_{l}} z}=\sum_{l=1}^{m^{(p)}} w_{l}(z) \vec{E}_{t, l} \tag{1}
\end{equation*}
$$

with

$$
\begin{equation*}
w_{l}(z)=a_{l} e^{-\jmath k_{z_{l}} z}+b_{l} e^{+\jmath k_{z_{l}} z}=\tilde{a}_{l}(z)+\tilde{b}_{l}(z) \tag{2}
\end{equation*}
$$

where $k_{z_{l}}$ is the propagation constant. We consider the application of (1) with (2) at a pair of neighboring cross-sectional planes $z_{p}$ and $z_{p+\Delta p} . m^{(p)}$ denotes the number of modes which have to be taken into account at the port $p$.
We get $\vec{E}_{t, l}\left(z_{p}\right)$ solving eigenvalue problems for the transmission lines. The scattering matrix $S$ (see [10]) is defined by

$$
\begin{equation*}
\overrightarrow{\vec{b}}_{\nu}=S \vec{a}_{\nu}, \quad \nu=1(1) m_{s} \tag{3}
\end{equation*}
$$

The dimension $m_{s}$ of this matrix is determined by the total number of modes at all ports. That means, we have to solve $m_{s}$ boundary value problems (see [9, 10]).

## 3 Boundary Value Problem

A three-dimensional boundary value problem can be formulated using the integral form of Maxwell's equations in the frequency domain [2] in order to compute the electromagnetic field:

$$
\begin{gather*}
\oint_{\partial \Omega} \vec{H} \cdot d \vec{s}=\jmath \omega \int_{\Omega}[\epsilon] \vec{E} \cdot d \vec{\Omega} \quad, \quad \oint_{\Omega}[\epsilon] \vec{E} \cdot d \vec{\Omega}=0 \\
\oint_{\partial \Omega} \vec{E} \cdot d \vec{s}=-\jmath \omega \int_{\Omega}[\mu] \vec{H} \cdot d \vec{\Omega} \quad, \quad \oint_{\Omega}[\mu] \vec{H} \cdot d \vec{\Omega}=0  \tag{4}\\
\vec{D}=[\epsilon] \vec{E}, \vec{B}=[\mu] \vec{H} \tag{5}
\end{gather*}
$$

The electric and magnetic flux densities $\vec{D}$ and $\vec{B}$ are complex functions of the spatial coordinates. $\omega=2 \pi f$ is the angular frequency of the sinusoidal excitation. $f$ denotes the frequency.

At the ports $p$ the transverse electric field $\vec{E}_{t}\left(z_{p}\right)$ is given by superposing weighted transmission line modes $\vec{E}_{t, l}\left(z_{p}\right)$ (see (1)):

$$
\begin{equation*}
\vec{E}_{t}\left(z_{p}\right)=\sum_{l=1}^{m^{(p)}} w_{l}\left(z_{p}\right) \vec{E}_{t, l}\left(z_{p}\right) \tag{6}
\end{equation*}
$$

All other parts of the surface of the computation domain are assumed to be an electric or a magnetic wall:

$$
\begin{equation*}
\vec{E} \times \vec{n}=0, \quad \vec{H} \times \vec{n}=0 \tag{7}
\end{equation*}
$$

We introduce a complex permittivity $[\epsilon]$ and a complex permeability $[\mu]$ diagonal tensor to obtain a reflection-free interface between the computational area and the lossy PML region:

$$
\begin{equation*}
[\epsilon]=(\epsilon)\left[\Lambda^{(\epsilon)}\right], \quad[\mu]=(\mu)\left[\Lambda^{(\mu)}\right] \tag{8}
\end{equation*}
$$

with

$$
\begin{equation*}
(\epsilon)=\operatorname{diag}\left(\epsilon_{x}, \epsilon_{y}, \epsilon_{z}\right), \quad(\mu)=\operatorname{diag}\left(\mu_{x}, \mu_{y}, \mu_{z}\right) . \tag{9}
\end{equation*}
$$

$\left[\Lambda^{(\epsilon)}\right]$ and $\left[\Lambda^{(\mu)}\right]$ are defined for a PML in $x$-, $y$-, or $z$-direction as follows $(\nu \in\{\epsilon, \mu\}$ ):

$$
\begin{gather*}
{\left[\Lambda^{(\nu)}\right]=\left\{\begin{array}{l}
{\left[\Lambda^{(\nu)}\right]_{x}=\operatorname{diag}\left(\frac{1}{\lambda_{\nu}}, \lambda_{\nu}, \lambda_{\nu}\right)} \\
{\left[\Lambda^{(\nu)}\right]_{y}=\operatorname{diag}\left(\lambda_{\nu}, \frac{1}{\lambda_{\nu}}, \lambda_{\nu}\right)} \\
{\left[\Lambda^{(\nu)}\right]_{z}=\operatorname{diag}\left(\lambda_{\nu}, \lambda_{\nu}, \frac{1}{\lambda_{\nu}}\right)}
\end{array}\right\} \text { with }}  \tag{10}\\
\lambda_{\nu}=1-\jmath \frac{\kappa_{\nu}}{\nu_{0} \omega} \text { and } \frac{\kappa_{\epsilon}}{\epsilon_{0}}=\frac{\kappa_{\mu}}{\mu_{0}} \tag{11}
\end{gather*}
$$

In case of overlapping at edges and corners the resulting PML tensor is the product of the PML tensors of the individual PML walls that form the edges and corners, respectively.

## 4 Maxwellian Grid Equations

Maxwellian grid equations are formulated for staggered nonequidistant rectangular grids (see Fig. 2) using the Finite Integration Technique with lowest order integration formulae [2, 14, 21]:

$$
\begin{equation*}
\oint_{\partial \Omega} \vec{f} \cdot d \vec{s} \rightarrow \sum\left( \pm f_{i} s_{i}\right), \quad \int_{\Omega} \vec{f} \cdot d \vec{\Omega} \rightarrow f \Omega, \quad \oint_{\Omega} \vec{f} \cdot d \vec{\Omega} \rightarrow \sum\left( \pm f_{i} \Omega_{i}\right) . \tag{12}
\end{equation*}
$$

The discretized form of (4) results in an equation for each field component. Presenting each equation using matrices provides a compact form:

$$
\begin{align*}
& \tilde{C} D_{\tilde{S} / / \tilde{\mu}} \vec{b}=\jmath \omega \epsilon_{0} \mu_{0} D_{\tilde{A} \epsilon} \vec{e}, \quad \tilde{S} D_{\tilde{A} \tilde{\epsilon}} \vec{e}=0,  \tag{13}\\
& C D_{s} \vec{e}=-\jmath \omega D_{A} \vec{b}, S D_{A} \vec{b}=0
\end{align*}
$$

The diagonal matrices $D_{\tilde{s} / \tilde{\mu}}, D_{\tilde{A} \tilde{\epsilon}}, D_{s}$, and $D_{A}$ represent all cell quantities. The so-called curl ( $C, \tilde{C}$ ) and source matrices $(S, \tilde{S})$ describe the topology of the two grids with the following properties (see [22]):

$$
\begin{equation*}
S C=0, \quad \tilde{S} \tilde{C}=0, \quad C=\tilde{C}^{T} \tag{14}
\end{equation*}
$$



Figure 2: Primary and dual grid

### 4.1 System of Linear Algebraic Equations

Using (14), eliminating the components of the magnetic flux density ( $\vec{b}$ ) in (13), and multiplying by $D_{s}^{1 / 2}$ yields a symmetric form of linear algebraic equations:

$$
\begin{equation*}
\left(D_{s}^{1 / 2} C^{T} D_{\tilde{s} / \tilde{\mu}} D_{A}^{-1} C D_{s}^{1 / 2}-k_{0}^{2} D_{\tilde{A} \tilde{\epsilon}}\right) D_{s}^{1 / 2} \vec{e}=0, \tag{15}
\end{equation*}
$$

where $k_{0}=\omega \sqrt{\epsilon_{0} \mu_{0}}$ denotes the wavenumber in vacuum. Moreover, the gradient of the electric field divergence

$$
\begin{equation*}
[\epsilon] \nabla\left([\epsilon]^{-2} \nabla \cdot[\epsilon] \vec{E}\right)=0 \tag{16}
\end{equation*}
$$

is equivalent to the matrix equation

$$
\begin{equation*}
\left(D_{s}^{-1 / 2} D_{\tilde{A} \tilde{\epsilon}} \tilde{S}^{T} D_{\tilde{V} \tilde{\epsilon} \tilde{S}}^{-1} \tilde{S} D_{\tilde{A} \tilde{\epsilon}} D_{s}^{-1 / 2}\right) D_{s}^{1 / 2} \vec{e}=0 \tag{17}
\end{equation*}
$$

The diagonal matrix $D_{\tilde{V} \tilde{\epsilon} \tilde{\epsilon}}$ is a volume matrix for the 8 partial volumes of the dual elementary cell.

The addition of Eqs. (15) and (17) yields the form

$$
\begin{equation*}
\left(B^{(G)}-k_{0}^{2} D^{(G)}\right) y^{(G)}=0 \tag{18}
\end{equation*}
$$

Taking into account the constitutive relations (5), the boundary conditions (7), and the transmission line modes $\vec{E}_{t, l}\left(z_{p}\right)$ (see (1)) we transform Eq. (18) into an inhomogeneous linear system of equations where its right-hand side depends on (6). For it, we use the notations given in Table 1. Thus, we get from Eq. (18):

$$
\binom{0}{B^{(G)}-k_{0}^{2} D^{(G)}}\binom{y_{E}^{(G)}}{y_{I}^{(G)}}=\left(\begin{array}{cc}
I_{E}-I_{E} & 0 \\
B_{E}^{(G)} & B_{I}^{(G)}-k_{0}^{2} D_{I}^{(G)}
\end{array}\right)\binom{y_{E}^{(G)}}{y_{I}^{(G)}}=0,
$$

Table 1: Notations

| $y^{(G)}=\left(y_{E}^{(G)}, y_{I}^{(G)}\right)^{T}$, vector of the unknown electric flux densities |  |  |
| :---: | :---: | :---: |
| $y_{E}^{(G)}=\left(y_{2 d}^{(G)}, y_{3 d}^{(G)}\right)^{T}$, given components of the solution vector $y^{(G)}$ |  |  |
| $y_{E}^{(G)}$ | external points, | $y_{2 d}^{(G)}$ solution of the |
|  | $\operatorname{dim}\left(y_{E}^{(G)}\right)=n_{E}$ | (G) $2 d$ eigenvalue problem |
| $y_{I}^{(G)}$ | internal points, | $y_{3 d}^{(G)}$ given boundary points |
|  | $\operatorname{dim}\left(y_{I}^{(G)}\right)=n_{I}, N=n_{E}+n_{I}$ | of the $3 d$ problem |
| $B^{(G)}$ | $=\left(B_{E}^{(G)}, 0\right)+\left(0, B_{I}^{(G)}\right)$ | $\operatorname{dim}\left(B^{(G)}\right)=\left(n_{I}, N\right)$, |
|  |  | $\operatorname{dim}\left(B_{E}^{(G)}\right)=\left(n_{I}, n_{E}\right),$ |
|  |  | $\operatorname{dim}\left(B_{I}^{(G)}\right)=\left(n_{I}, n_{I}\right)$ |
| $D^{(G)}$ | $=\left(0, D_{I}^{(G)}\right)$ | $\operatorname{dim}\left(D^{(G)}\right)=\left(n_{I}, N\right)$, |
|  |  | $\operatorname{dim}\left(D_{I}^{(G)}\right)=\left(n_{I}, n_{I}\right)$ |
| $I_{E}$ | identity | $\operatorname{dim}\left(I_{E}\right)=\left(n_{E}, n_{E}\right)$ |

$$
\begin{gather*}
\left(\begin{array}{cc}
I_{E} & 0 \\
0 & B_{I}^{(G)}-k_{0}^{2} D_{I}^{(G)}
\end{array}\right)\binom{y_{E}^{(G)}}{y_{I}^{(G)}}+\left(\begin{array}{cc}
-I_{E} & 0 \\
B_{E}^{(G)} & 0
\end{array}\right)\binom{y_{E}^{(G)}}{y_{I}^{(G)}}=0, \\
\left(\begin{array}{cc}
I_{E} & 0 \\
0 & B_{I}^{(G)}-k_{0}^{2} D_{I}^{(G)}
\end{array}\right)\binom{y_{E}^{(G)}}{y_{I}^{(G)}}=\binom{y_{E}^{(G)}}{-B_{E}^{(G)} y_{E}^{(G)}}=\binom{y_{E}^{(G)}}{c_{I}^{(G)}}=c^{(G)} . \tag{19}
\end{gather*}
$$

Therefore, the systems of linear algebraic equations

$$
\begin{equation*}
\left(B_{I}^{(G)}-k_{0}^{2} D_{I}^{(G)}\right) y_{I}^{(G)}=c_{I}^{(G)} \tag{20}
\end{equation*}
$$

are to be solved. The number of such systems (20) is $m_{s}$ (see Section 2).

### 4.2 Eigenvalue Problem

The vector $y_{2 d}^{(G)}$ is the solution of the $2 d$ eigenvalue problem. In the following, we consider a longitudinally homogeneous transmission line. Thus, any field can be expanded into a sum of so-called modal fields which vary exponentially in the longitudinal direction:

$$
\begin{equation*}
\vec{E}(x, y, z \pm 2 h)=\vec{E}(x, y, z) e^{\mp j k_{z} 2 h} \tag{21}
\end{equation*}
$$

where $k_{z}$ is the propagation constant, and $2 h$ is the length of an elementary cell in $z$-direction (see Fig. 3). Thus, we get a two-dimensional eigenvalue problem for the transverse electric fields $\vec{y}=\vec{E}_{t, l}\left(z_{p}\right), \quad l=1(1) m^{(p)}$, (see (6)) on the transmission line region:

$$
\begin{equation*}
A \vec{y}=\gamma \vec{y}, \gamma=e^{-\jmath k_{z} 2 h}+e^{+\jmath k_{z} 2 h}-2=-4 \sin ^{2}\left(h k_{z}\right) . \tag{22}
\end{equation*}
$$

A detailed derivation of the eigenvalue problem can be found in [11, 12, 13, 15].


Figure 3: Transmission line

## 5 Symmetric Band Lanczos Process

We consider the iterative solution of large systems of linear algebraic equations which not only have multiple right-hand sides, but also have multiple shifts for each right-hand side. The generalized form of Eq. (20) is the problem

$$
\begin{equation*}
\left(\alpha_{j} B+\beta_{j} D\right) y^{(j, k)}=c^{(k)}, \quad \alpha_{j}, \beta_{j} \in \mathbb{C}, \tag{23}
\end{equation*}
$$

with $j=1, \ldots, n_{s}$ and $k=1, \ldots, n_{b}$. Let $n_{b}$ be the number of right-hand sides and $n_{s}$ be the number of shifts. The Matrix $B$ is complex, symmetric, and indefinite. A standard way to solve systems with multiple right-hand sides is to use a block approach $[18,8,3]$.

Eq. (23) can be transformed into

$$
\begin{equation*}
\left(D^{-1 / 2} B D^{-1 / 2}+\alpha_{j}^{-1} \beta_{j} I\right)\left(\alpha_{j} D^{1 / 2} y^{(j, k)}\right)=D^{-1 / 2} c^{(k)} . \tag{24}
\end{equation*}
$$

Thus, we get then the common equation

$$
\begin{equation*}
A^{(j)} x^{(j, k)}=\left(A+\sigma_{j} I\right) x^{(j, k)}=b^{(k)}, \quad \sigma_{j} \in \mathbb{C} \tag{25}
\end{equation*}
$$

i.e.,

$$
\left(\begin{array}{ccccc}
A^{(1)} & & & & \\
& \ddots & & 0 & \\
& & A^{(j)} & & \\
& 0 & & \ddots & \\
& & & & A^{\left(n_{s}\right)}
\end{array}\right)\left(\begin{array}{ccccc}
x^{(1,1)} & \cdots & x^{(1, k)} & \cdots & x^{\left(1, n_{b}\right)} \\
\vdots & & \vdots & & \vdots \\
x^{(j, 1)} & \cdots & x^{(j, k)} & \cdots & x^{\left(j, n_{b}\right)} \\
\vdots & & \vdots & & \vdots \\
x^{\left(n_{s}, 1\right)} & \cdots & x^{\left(n_{s}, k\right)} & \cdots & x^{\left(n_{s}, n_{b}\right)}
\end{array}\right)=
$$

$$
\left(\begin{array}{ccccc}
b^{(1)} & \cdots & b^{(k)} & \cdots & b^{\left(n_{b}\right)} \\
\vdots & & \vdots & & \vdots \\
b^{(1)} & \cdots & b^{(k)} & \cdots & b^{\left.n_{b}\right)} \\
\vdots & & \vdots & & \vdots \\
b^{(1)} & \cdots & b^{(k)} & \cdots & b^{\left(n_{b}\right)}
\end{array}\right) \quad \begin{gathered}
1 \\
j \\
\\
n_{s}
\end{gathered}
$$

with

$$
A=D^{-1 / 2} B D^{-1 / 2}, x^{(j, k)}=\alpha_{j} D^{1 / 2} y^{(j, k)}, b^{(k)}=D^{-1 / 2} c^{(k)}, \text { and } \sigma_{j}=\alpha_{j}^{-1} \beta_{j} .
$$

For the special case that all right-hand sides in (25) are identical, i.e.,

$$
\begin{equation*}
x^{(j)}=x^{(j, k)}, b^{(k)}=b \text { for all } k=1, \ldots, n_{b}, \tag{26}
\end{equation*}
$$

it is straightforward to exploit the shift structure when solving the $n_{s}$ systems by Krylov subspace methods. We use the initial guess $x_{0}^{(j)}=0$ for all $j$. In this case, the Krylov subspaces for all $n_{s}$ systems are identical:

$$
\begin{equation*}
\mathcal{K}_{m}\left(A+\sigma_{j} I, b\right)=\mathcal{K}_{m}(A, b) \text { for all } j=1, \ldots, n_{s} \text { and } m \geq 1 \tag{27}
\end{equation*}
$$

This means that the computation of suitable basis vectors for the underlying Krylov subspaces has to be performed only once. A lot of Krylov subspace methods have been developed for shifted matrix problems. A detailed derivation of a shifted coupled two-term algorithm without look-ahead for

$$
\begin{equation*}
A^{(j)} x^{(j)}=\left(A+\sigma_{j} I\right) x^{(j)}=b, \quad \sigma_{j} \in \mathbb{C}, \quad j=1, \ldots, n_{s}, \tag{28}
\end{equation*}
$$

can be found in [19].

### 5.1 Symmetric Block-QMR Method

We first describe the general form of block Krylov subspace iterative methods for the simultaneous solution of linear systems with multiple right-hand sides (see [8, 16]). Solving the $n_{b}$ systems

$$
\begin{equation*}
A x^{(j)}=b^{(j)}, j=1, \ldots, n_{b}, \tag{29}
\end{equation*}
$$

where $A \in \mathbb{C}^{N \times N}, x^{(j)} \in \mathbb{C}^{N}$, and $b^{(j)} \in \mathbb{C}^{N}$ is equivalent to solving the block system of linear equations

$$
\begin{equation*}
A X=B \tag{30}
\end{equation*}
$$

where $B=\left[b^{(1)} b^{(2)} \ldots b^{\left(n_{b}\right)}\right] \in \mathbb{C}^{N \times n_{b}}$ and $X=\left[x^{(1)} x^{(2)} \ldots x^{\left(n_{b}\right)}\right] \in \mathbb{C}^{N \times n_{b}}$.
The block Krylov method generates a sequence of block iterates

$$
\begin{equation*}
X_{\mu}=\left[x_{\mu}^{(1)} x_{\mu}^{(2)} \ldots x_{\mu}^{\left(n_{b}\right)}\right], \mu=1,2, \ldots, \tag{31}
\end{equation*}
$$

where $x_{\mu}^{(j)} \in x_{0}^{(j)}+\mathcal{K}_{\mu}(A, R)$ for each $j=1, \ldots, n_{b}$ and $R \in \mathbb{C}^{N \times n_{b}}$.
Here, $R=R_{0}=B-A X_{0}$ is the block of initial residual vectors and $\mathcal{K}_{\mu}(A, R)$ denotes the
$\mu$-th block Krylov subspace generated by $A$ and $R$, i.e., the first $\mu$ linearly independent vectors in the block Krylov sequence

$$
\begin{equation*}
R, A R, \ldots, A^{k-1} R, \ldots \tag{32}
\end{equation*}
$$

span the subspace $\mathcal{K}_{\mu}(A, R)$.
By scanning the vectors in (32) from left to right and deleting each vector that is either linearly dependent or almost linearly dependent on previous vectors, we obtain the so-called deflated block Krylov sequence

$$
\begin{equation*}
R_{1}, A R_{2}, \ldots, A^{k-1} R_{k}, \ldots \tag{33}
\end{equation*}
$$

For $k=1$, we set $R=R_{0}$.
In addition to the matrices $A$ and $R$, the Lanczos-type process requires a second sequence of vectors that span the block Krylov subspace $\mathcal{K}_{\mu}\left(A^{T}, L\right)$ generated by $A^{T}$ and $L$, i.e., the Lanczos-type process has produced right and left Lanczos vectors

$$
\begin{equation*}
\left\{v_{1}, \ldots, v_{\mu}\right\} \quad \text { and } \quad\left\{w_{1}, \ldots, w_{\mu}\right\} \tag{34}
\end{equation*}
$$

respectively, with

$$
\mathcal{K}_{\mu}(A, R)=\operatorname{span}\left\{v_{1}, \ldots, v_{\mu}\right\}
$$

and

$$
\mathcal{K}_{\mu}\left(A^{T}, L\right)=\operatorname{span}\left\{w_{1}, \ldots, w_{\mu}\right\}
$$

The vectors (34) are constructed to be biorthogonal, i.e.,

$$
w_{i}^{T} v_{j}=\left\{\begin{array}{ll}
\delta_{j} \neq 0 & \text { if } \quad i=j \\
0 & \text { if } \quad i \neq j
\end{array} \quad \text { for all } \quad i, j=1, \ldots, \mu\right.
$$

An important feature of the Lanczos-type process is that the vectors (34) are generated by means $\left(2 n_{b}+1\right)$-term recurrences. These recurrences can be summarized in compact form as follows:

$$
\begin{equation*}
A V_{\mu}=V_{n} T_{\mu}+V_{\mu}^{d f}, \mu=n-m_{c r} \geq 1 \tag{35}
\end{equation*}
$$

Here, $m_{c r}$ denotes the reduced size of the current block in the deflated block Krylov sequence (33) and $n$ is the iteration index. $V_{n}=\left[\begin{array}{llll}v_{1} & v_{2} & \ldots & v_{n}\end{array}\right] \in \mathbb{C}^{N \times n}$ contains the first $n$ basis vectors, $T_{\mu} \in \mathbb{C}^{n \times \mu}$ is a banded matrix with lower and upper band width $n_{b}+1$ containing the recurrence coefficients, and the matrix $V_{\mu}^{d f} \in \mathbb{C}^{N \times \mu}$ consists of mostly zero column vectors and a few nonzero columns with very small entries.
Each deflation step in the block Krylov sequence (33) reduces the lower band width of $T_{\mu}$ by one. Additionally, the recurrences (35) are complemented by recurrences for the initial block $R$, i.e.,

$$
\begin{equation*}
V_{m_{1}} \Theta+V_{0}^{d f}=R, \quad V_{m_{1}} \in \mathbb{C}^{N \times m_{1}}, \quad m_{1} \leq \operatorname{rank}(R) \leq n_{b} \tag{36}
\end{equation*}
$$

The matrix $\Theta \in \mathbb{C}^{m_{1} \times n_{b}}$ contains recurrence coefficients up to iteration when $\mu=0$. The matrix $V_{0}^{d f} \in \mathbb{C}^{N \times n_{b}}$ consists of zero vectors and deflated vectors.
Notice that (35) only holds for $\mu \geq 1$, or equivalent for $n>m_{1}$, where $m_{1}$ is the size of the first block in the sequence (33), i.e., $m_{1}=\operatorname{rank}\left(R_{1}\right)$.

The left Lanczos vectors in (34) are generated by means of recurrences similar to (35) and (36). All possible block iterates (31) can be represented in the form

$$
\begin{equation*}
X_{\mu}=X_{0}+V_{\mu} Z_{\mu}, \quad Z_{\mu} \in \mathbb{C}^{\mu \times n_{b}} \tag{37}
\end{equation*}
$$

After each deflation in the $v$ sequence, the size of the current block in the corresponding deflated block Krylov sequence (33) decrease by one. Additionally, one of the linear system of (30) can be dropped from subsequent block QMR iteration (see [8]). The solution of all such deflated systems can be recovered from the solution of the converged remaining linear systems. Therefore, one of the columns from the matrices $X_{\mu}, X_{0}$, and $Z_{\mu}$ in (37) can be dropped and only the remaining iterates need to be updated. Then, for the current block iterate $X_{\mu}^{c r}$, we can rewrite (37) as

$$
\begin{equation*}
X_{\mu}^{c r}=X_{0}^{c r}+V_{\mu} Z_{\mu}, \quad Z_{\mu} \in \mathbb{C}^{\mu \times m_{c r}}, \quad m_{c r} \leq n_{n} \tag{38}
\end{equation*}
$$

The residual block $R_{\mu}^{c r}$ corresponding to $X_{\mu}^{c r}$ is given by

$$
\begin{align*}
R_{\mu}^{c r} & =B^{c r}-A X_{\mu}^{c r} \\
& =R_{0}^{c r}-V_{n} T_{\mu} Z_{\mu}-V_{\mu}^{d f} Z_{\mu} \\
& =V_{n}\left(\left[\begin{array}{l}
\Theta^{c r} \\
0
\end{array}\right]-T_{\mu} Z_{\mu}\right)-V_{\mu}^{d f} Z_{\mu} \tag{39}
\end{align*}
$$

The matrix $\Theta^{c r} \in \mathbb{C}^{m_{1} \times m_{c r}}$ consists of those columns of $\Theta$ in (36) that correspond to columns of $X_{\mu}$ retained in $X_{\mu}^{c r}$.
The free parameter matrix $Z_{\mu}$ in (38) is determined as the solution of the matrix least-squares problem

$$
\left\|\left[\begin{array}{l}
\Theta^{c r}  \tag{40}\\
0
\end{array}\right]-T_{\mu} Z_{\mu}\right\|=\min _{Z \in \mathbb{C}^{\mu \times m_{c r}}}\left\|\left[\begin{array}{l}
\Theta^{c r} \\
0
\end{array}\right]-T_{\mu} Z\right\|
$$

in the Euclidean norm. The matrix least-squares problem (40) is solved by means standard techniques based on QR factorization of $T_{\mu}$. This allows to obtain the solution $Z_{\mu}$ by updating the solution $Z_{\mu-1}$ from the previous step. Implementation details are described in [8].
The banded matrix $T_{\mu}$ is factorized into a unitary matrix $\left(Q_{\mu}\right)^{H}$ and a nonsingular upper triangular matrix $U_{\mu}$ :

$$
T_{\mu}=\left(Q_{\mu}\right)^{H}\left[\begin{array}{c}
U_{\mu} \\
0
\end{array}\right], \quad Q_{\mu} \in \mathbb{C}^{n \times n}, \quad U_{\mu} \in \mathbb{C}^{\mu \times \mu}
$$

The lower bandwidth of the matrix $T_{\mu}$ is at most $n_{b}+1$. Its QR factorization can be updated using $n_{b}$ Givens rotations at each iteration of the block QMR method. The update of the block QMR iterates can be performed by means the following quantities:

$$
\begin{gathered}
{\left[\begin{array}{l}
t_{\mu}^{c r} \\
\tau_{\mu}^{c r}
\end{array}\right]=Q_{\mu}\left[\begin{array}{l}
\Theta^{c r} \\
0
\end{array}\right], \quad t_{\mu}^{c r} \in \mathbb{C}^{\mu \times m_{c r}}, \quad \tau_{\mu}^{c r} \in \mathbb{C}^{m_{c r} \times m_{c r}},} \\
t_{\mu}^{c r}=\left[\begin{array}{l}
t_{\mu-1}^{c r} \\
y_{\mu}^{T}
\end{array}\right], \quad y_{\mu}^{T} \in \mathbb{C}^{1 \times m_{c r}} .
\end{gathered}
$$

The solution $Z_{\mu}$ of (40) is given by

$$
Z_{\mu}=U_{\mu}^{-1} t_{\mu}^{c r}
$$

We also define a sequence of direction vectors $p_{1}, p_{2}, \ldots, p_{\mu}, \mu=1,2, \ldots$, with

$$
P_{\mu}=\left[\begin{array}{llll}
p_{1} & p_{2} & \ldots & p_{\mu}
\end{array}\right]=V_{\mu} U_{\mu}^{-1} .
$$

Finally, the block QMR iterates can be updated by means

$$
\begin{aligned}
X_{\mu}^{c r} & =X_{0}^{c r}+V_{\mu} U_{\mu}^{-1} t_{\mu}^{c r} \\
& =X_{\mu-1}^{c r}+p_{\mu} y_{\mu}^{T} .
\end{aligned}
$$

It is well known that the classical Lanczos process for general matrices $A$ simplifies when applied to complex symmetric matrices. The resulting Lanczos process only involves one sequence of Lanczos vectors. The resulting simplified block QMR method does not require matrixvector multiplications with $A^{T}$.

Now, we consider an adaption of block QMR to shifted systems, i.e., systems of the type

$$
\begin{equation*}
\left(A+\sigma_{j} I\right) x^{(j)}=b^{(j)}, \quad \sigma_{j} \in \mathbb{C}, \quad j=1, \ldots, n_{b} . \tag{41}
\end{equation*}
$$

We sketch the general case (41) or the case (28) if non zero initial guesses, $x_{0}^{(j)} \neq 0$, are chosen. Then, the Krylov subspaces depend on $j$. We use the Lanczos process to still exploit the shift structure of (41). Using the right initial block

$$
\begin{equation*}
R=R_{0}=\left[r_{0}^{(1)} r_{0}^{(2)} \ldots r_{0}^{\left(n_{b}\right)}\right], \quad r_{0}^{(j)}=b^{(j)}-\left(A+\sigma_{j} I\right) x_{0}^{(j)}, \tag{42}
\end{equation*}
$$

we generate basis vectors $\left\{\begin{array}{llll}v_{1} & v_{2} & \ldots & v_{\mu}\end{array}\right\}$ for th $\mu$-th $\operatorname{Krylov}$ subspace $\mathcal{K}_{\mu}(A, R)$. The $\mu$-th iterate $x_{\mu}^{(j)}$ for the $j$-th system in (41) is then defined by

$$
\begin{equation*}
x_{\mu}^{(j)}=x_{0}^{(j)}+V_{\mu} z_{\mu}^{(j)}, \tag{43}
\end{equation*}
$$

where $z_{\mu}^{(j)} \in \mathbb{C}^{\mu}$ is the solution of the least-squares problem

$$
\left\|\left[\begin{array}{l}
\Theta_{j}^{c r}  \tag{44}\\
0
\end{array}\right]-T_{\mu}^{(j)} z_{\mu}^{(j)}\right\|=\min _{z \in \mathbb{C}^{\mu}}\left\|\left[\begin{array}{l}
\Theta_{j}^{c r} \\
0
\end{array}\right]-T_{\mu}^{(j)} z\right\| .
$$

Using (35) and $A^{(j)}=A+\sigma_{j} I$, we get

$$
\begin{equation*}
A^{(j)} V_{\mu}=V_{n} T_{\mu}^{(j)}+V_{\mu}^{d f} \tag{45}
\end{equation*}
$$

Furthermore,

$$
\begin{aligned}
\left(A+\sigma_{j} I\right) V_{\mu} & =A V_{\mu}+\sigma_{j} V_{\mu} \\
& =V_{n} T_{\mu}+\sigma_{j} V_{\mu}+V_{\mu}^{d f} \\
& =V_{n}\left(T_{\mu}+\sigma_{j}\left[\begin{array}{c}
I_{\mu} \\
0
\end{array}\right]\right)+V_{\mu}^{d f} .
\end{aligned}
$$

With $T_{\mu}^{(j)}=T_{\mu}+\sigma_{j}\left[\begin{array}{c}I_{\mu} \\ 0\end{array}\right]$ we get from (44) the least-squares problem

$$
z_{\mu}^{(j)}=\arg \min _{z \in \mathbb{C}^{\mu}} \|\left[\begin{array}{l}
\Theta_{j}^{c r}  \tag{46}\\
0
\end{array}\right]-\left(\left(T_{\mu}+\sigma_{j}\left[\begin{array}{c}
I_{\mu} \\
0
\end{array}\right]\right) z \|\right.
$$

$\operatorname{In}(46), T_{\mu}$ is the matrix of Lanczos recurrence coefficients given by (35). The vector $\Theta_{j}^{c r}$ is defined as the $j$-th column of the matrix $\Theta^{c r}$ (see (36) and (39)).

### 5.2 Symmetric Band Lanczos Method

The symmetric band Lanczos process is an extension of the classical Lanczos algorithm for symmetric matrices and single starting vectors to multiple starting vectors (see [1, 4]). First, we consider the $n_{b}$ linear systems (29), where $A$ is a complex symmetric matrix and $R=$ $\left[r^{(1)} r^{(2)} \ldots r^{\left(n_{b}\right)}\right]$ is a block of $n_{b}$ right starting vectors. The symmetric band Lanczos process generates orthogonal basis vectors for the subspaces spanned by leading columns of the block Krylov matrix (see (32)):

$$
\begin{equation*}
\mathcal{K}(A, R)=\left[R A R \ldots A^{N-1} R\right] . \tag{47}
\end{equation*}
$$

The columns of the matrix $\mathcal{K}(A, R)$ in (47) are all vectors of length $N$, and thus at most $N$ of them are linearly independent. One needs to perform so-called deflation of linearly dependent or in some sence almost linearly dependent vectors.

After $n$ iterations, the algorithm has generated the first $n$ Lanczos vectors

$$
\begin{equation*}
v_{1}, v_{2}, \ldots, v_{n} \in \mathbb{C}^{N} \tag{48}
\end{equation*}
$$

It will be convenient to introduce the notation

$$
V_{n}=\left[\begin{array}{llll}
v_{1} & v_{2} & \ldots & v_{n} \tag{49}
\end{array}\right], \quad V_{n}^{T} V_{n}=I_{n},
$$

for the matrices whose columns are just the right Lanczos vectors (48). In addition to (48), the algorithm has constructed the $m_{c r}=m_{c r}(n)$ vectors

$$
\begin{equation*}
\hat{v}_{n+1}, \hat{v}_{n+2}, \ldots, \hat{v}_{n+m_{c r}} \in \mathbb{C}^{N} \tag{50}
\end{equation*}
$$

that are candidates for the next $m_{c r}$ Lanczos vectors, $v_{n+1}, v_{n+2}, \ldots, v_{n+m_{c r}}$. Here $m_{c r}$ is the current block size. At the initialization phase, i.e., $n=0$, we have

$$
m_{c r}=n_{b} \quad \text { and } \quad \hat{v}_{j}=r^{(j)}, 1 \leq j \leq n_{b} .
$$

Within the algorithm, $m_{c r}$ is reduced by one every time a deflation occurs. The vectors (50) are constructed so that they satisfy the orthogonality relations

$$
\begin{equation*}
V_{n}^{T} \hat{v}_{n+j}=0 \text { for all } j=1, \ldots, m_{c r} \tag{51}
\end{equation*}
$$

The recurrences can be summarized compactly as follows:

$$
\begin{equation*}
A V_{n}=V_{n} T_{n}+[\underbrace{0 \ldots 0}_{n-m_{c r}} \underbrace{\hat{v}_{n+1} \hat{v}_{n+2} \ldots \hat{v}_{n+m_{c r}}}_{m_{c r}}]+\hat{V}_{n}^{d f} \tag{52}
\end{equation*}
$$

The entries of the matrix $T_{n}, T_{n} \in \mathbb{C}^{n \times n}$, are chosen so that the orthogonality conditions (49) and (51) are satisfied. The matrix $\hat{V}_{n}^{d f}$ contains mostly zero columns together with the $\hat{v}_{j}$ vectors that have been deflated during the first $n$ iterations. We remark that $n_{b}-m_{c r}$ is the number of deflated $\hat{v}_{j}$ vectors. It turns out that orthogonality only has to be explicitly enforced among $2 m_{c r}+1$ consecutive Lanczos vectors and, once deflation has occurred, against $n_{b}-m_{c r}$ fixed earlier right Lanczos vectors. Thus, the matrix $T_{n}$ is essentially banded. More precisely, $T_{n}$ has
lower as well as upper bandwidth $m_{c r}+1$. Additionally, the recurrences (52) are complemented by recurrences for the initial block $R$, i.e.,

$$
\begin{equation*}
R=R_{0}=V_{m_{1}} \Theta+\hat{V}_{0}^{d f}, \quad V_{m_{1}} \in \mathbb{C}^{N \times m_{1}}, \quad m_{1} \leq \operatorname{rank}(R) \leq n_{b} \tag{53}
\end{equation*}
$$

Here, $m_{1} \leq n_{b}$ denotes the number of columns of the block $R$ that have not been deflated. The matrix $\hat{V}_{0}^{d f} \in \mathbb{C}^{N \times n_{b}}$ contains the $n_{b}-m_{1}$ deflated starting vectors and $m_{1}$ zero vectors as columns, and $\Theta \in \mathbb{C}^{m_{1} \times n_{b}}$ is an upper triangular matrix whose entries are chosen such that the columns of $V_{m_{1}}$ are orthonormal.
If no deflation has occured or if only exact deflation is performed,i.e.,

$$
\hat{V}_{0}^{d f}=0 \quad \text { and } \quad \hat{V}_{n}^{d f}=0
$$

then

$$
\begin{equation*}
V_{n}^{T} A V_{n}=V_{n}^{T} V_{n} T_{n}=T_{n} \tag{54}
\end{equation*}
$$

Furthermore, the so-called $n$-th Lanczos matrix $T_{n}$ is symmetric:

$$
\begin{equation*}
T_{n}^{T}=\left(V_{n}^{T} A V_{n}\right)^{T}=V_{n}^{T} A V_{n}=T_{n} \tag{55}
\end{equation*}
$$

In the case $\hat{V}_{n}^{d f} \neq 0$, we get

$$
\tilde{T}_{n}=V_{n}^{T} A V_{n}+V_{n}^{T} \hat{V}_{n}^{d f}=T_{n}+V_{n}^{T} \hat{V}_{n}^{d f}, \quad \tilde{T}_{n} \in \mathbb{C}^{n \times n}
$$

The matrix $\tilde{T}_{n}$ consists of a symmetric banded part, $T_{n}$, with decreasing bandwidth and a spiked part with potentially nonzero elements only in rows and columns with index $i \in \mathcal{I}$ and outside of the banded part. The index set $\mathcal{I}$ records the positions of the potentially nonzero columns of $\hat{V}_{n}^{d f}$ due to deflation.
The matrix $T_{n}$ is decomposed in $T_{n}=L_{n} U_{n}$, where $L_{n} \in \mathbb{C}^{n \times n}$ and $U_{n} \in \mathbb{C}^{n \times n}$ are lower and upper triangular matrices,respectively. The unit upper triangular $U_{n}$ consists of a banded part with bandwidth decreasing from $n_{b}+1$ to $m_{c r}+1$. Analogous, the lower triangular $L_{n}$ consists of a banded part with bandwidth decreasing from $n_{b}+1$ to $m_{c r}+1$ :

$$
L_{n}=\left(\begin{array}{ccc}
l_{11} & & 0  \tag{56}\\
\vdots & \ddots & \\
l_{1 n} & \cdots & l_{n n}
\end{array}\right) \quad, \quad U_{n}=\left(\begin{array}{ccc}
1 & \cdots & u_{1 n} \\
& \ddots & \vdots \\
0 & & 1
\end{array}\right)
$$

where $l_{i j}=u_{j i}=0$ for $j+m_{c r}<i \leq n, j=1, \ldots, n-m_{c r}$, and $n>m_{c r}$.
We generate a second set of vectors

$$
\begin{equation*}
p_{1}, p_{2}, \ldots, p_{n} \in \mathbb{C}^{N}, \quad P_{n}=\left[p_{1} p_{2} \ldots p_{n}\right] \tag{57}
\end{equation*}
$$

that span the same subspaces as $V_{n}$, i.e.,

$$
\operatorname{span}\left\{v_{1}, v_{2}, \ldots, v_{j}\right\}=\operatorname{span}\left\{p_{1}, p_{2}, \ldots, p_{j}\right\} \quad \text { for all } \quad 1 \leq j \leq n
$$

Then we set

$$
\begin{equation*}
V_{n}=P_{n} U_{n} \tag{58}
\end{equation*}
$$

All block iterates can be represented in the form

$$
\begin{equation*}
X_{n}=X_{0}+V_{n} Z_{n}, \quad Z_{n} \in \mathbb{C}^{n \times n_{b}} \tag{59}
\end{equation*}
$$

After each deflation, for the current block iterate $X_{n}^{c r}$, we can rewrite (59) as

$$
\begin{align*}
X_{n}^{c r} & =X_{0}^{c r}+V_{n} Z_{n}, \quad Z_{n} \in \mathbb{C}^{n \times m_{c r}} \\
& =X_{0}^{c r}+P_{n} U_{n} Z_{n} \\
& =X_{0}^{c r}+P_{n} Y_{n}, \quad Y_{n}=U_{n} Z_{n} \in \mathbb{C}^{n \times m_{c r}} \\
& =X_{n-1}^{c r}+p_{n}\left(e_{n}^{T} Y_{n}\right) \tag{60}
\end{align*}
$$

Using (52), (53), (58) - (60), and $\hat{V}_{n}^{c r}=\left[\begin{array}{llllll}0 & \ldots & 0 & \hat{v}_{n+1} & \ldots & \hat{v}_{n+m_{c r}}\end{array}\right]$ the residual block $R_{n}^{c r}$ corresponding to $X_{n}^{c r}$ can be represented as follows:

$$
\begin{align*}
R_{n}^{c r} & =B^{c r}-A X_{n}^{c r} \\
& =B^{c r}-A X_{0}^{c r}-A V_{n} Z_{n} \\
& =R_{0}^{c r}-V_{n} T_{n} Z_{n}-\hat{V}_{n}^{c r} Z_{n}-\hat{V}_{n}^{d f} Z_{n} \\
& =V_{m_{c r}} \Theta^{c r}-V_{n} L_{n} Y_{n}-\hat{V}_{n}^{c r} Z_{n}-\hat{V}_{n}^{d f} Z_{n} \\
& =V_{n}\left(\left[\begin{array}{l}
\Theta^{c r} \\
0
\end{array}\right]-L_{n} Y_{n}\right)-\hat{V}_{n}^{c r} Z_{n}-\hat{V}_{n}^{d f} Z_{n} . \tag{61}
\end{align*}
$$

Using (61), $Y_{n}$ is the unique solution of the least-squares problem

$$
Y_{n}=\arg \min _{Y \in \mathbb{C}^{n \times m} m_{c r}}\left\|\left[\begin{array}{l}
\Theta^{c r}  \tag{62}\\
0
\end{array}\right]-L_{n} Y\right\|
$$

We describe now the basic steps of the QMR approach. We first consider the solution of the least-squares problem (62) for the case when no deflation occurs at iteration $n$, while generating $v_{n}$ :

$$
V_{n}=\left[\begin{array}{ll}
V_{n-1} & v_{n}
\end{array}\right]
$$

Using (56), we consider the decomposition of the $n$-th Lanczos matrix $T_{n}$ :

$$
\left.\begin{gathered}
T_{n}=\left(\begin{array}{ccc|c} 
& & & t_{1 n} \\
& T_{n-1} & & \vdots \\
& & & t_{n-1 n}
\end{array}\right), \quad t_{i j}=t_{j i}, \quad 1 \leq i, j \leq n, \\
\hline t_{n 1} \\
\cdots
\end{gathered} t_{n n-1} \right\rvert\, t_{n n}, ~\left(\begin{array}{cc|c} 
& L_{n-1} & 0 \\
\hline l_{n 1} & \cdots & l_{n n-1}
\end{array}\right)\left(\begin{array}{c|c}
l_{n n}
\end{array}\right)\left(\begin{array}{c}
u_{1 n} \\
U_{n-1} \\
\hline 0 \\
\hline
\end{array}\right) .
$$

Due to the fact that the lower and upper triangular $L_{n}$ and $U_{n}$, respectively, consist of a banded
part with bandwidth $m_{c r}+1$ we get the following terms:

$$
\begin{aligned}
& l_{n i}=t_{n i}-\sum_{j=j^{*}}^{i-1} l_{n j} u_{j i}, \quad u_{i i}=1 \\
& u_{i n}=\frac{1}{l_{i i}}\left(t_{i n}-\sum_{j=j^{*}}^{i-1} l_{i j} u_{j n}\right) \\
& u_{n n}=1 \\
& l_{n n}=t_{n n}-\sum_{i=i^{*}}^{n-1} l_{n i} u_{i n}
\end{aligned}
$$

for $i^{*}=\max \left(1, n-m_{c r}\right), j^{*}=\max \left(1, n-m_{c r}\right)$, and $i=i^{*}, \ldots, n-1$. For the solution of the least-squares problem (62) we get

$$
\begin{aligned}
& L_{n} Y_{n}=S_{n}, \quad L_{n} \in \mathbb{C}^{n \times n}, \quad Y_{n} \in \mathbb{C}^{n \times m_{c r}}, \quad S_{n} \in \mathbb{C}^{n \times m_{c r}}, \text { i.e. }, \\
& \left(\right)\left(\begin{array}{ccc}
\frac{Y_{n-1}}{} \\
\hline y_{n 1} & \cdots & y_{n m_{c r}}
\end{array}\right)=\left(\begin{array}{ccc} 
& S_{n-1} \\
\hline s_{n 1} & \cdots & s_{n m_{c r}}
\end{array}\right),
\end{aligned}
$$

and

$$
\begin{align*}
L_{n-1} Y_{n-1} & =S_{n-1} \\
y_{n i} & =\frac{1}{l_{n n}}\left(s_{n i}-\left(l_{n j^{*}} \cdots l_{n n-1}\right) Y_{n-1}^{(i)}\right), \quad i=1, \ldots, m_{c r} \tag{63}
\end{align*}
$$

where $Y_{n-1}^{(i)}$ is the $i$-th column of the matrix $Y_{n-1}$. We set $S_{n}=\Theta^{c r}$ for $n=m_{c r}$ and $s_{n i}=0$ for $i=1, \ldots, m_{c r}$ and $n>m_{c r}$.

Next, we describe the effect of deflation in the Lanczos-type algorithm for updating the block iterates $X_{n}^{c r}$. In this case, we set $v_{n}=0$, which implies that

$$
V_{n}=\left[\begin{array}{ll}
V_{n-1} & 0
\end{array}\right] .
$$

$V_{n}$ now only has column rank $n-1$.

$$
\begin{gathered}
T_{n}=V_{n}^{T} A V_{n}=\left[\begin{array}{c}
V_{n-1}^{T} \\
0
\end{array}\right] A\left[V_{n-1} 0\right]=\left(\begin{array}{c|c}
T_{n-1} & 0 \\
\hline 0 & 0
\end{array}\right) \\
T_{n}=L_{n} U_{n}=\left(\begin{array}{c|c}
L_{n-1} U_{n-1} & 0 \\
\hline 0 & 0
\end{array}\right)=\left(\begin{array}{c|c}
L_{n-1} & 0 \\
\hline 0 & 0
\end{array}\right)\left(\begin{array}{c|c}
U_{n-1} & 0 \\
\hline 0 & 1
\end{array}\right)
\end{gathered}
$$

For the case when no deflation occurs we get

$$
L_{n} Y_{n}=\left[\begin{array}{l}
t_{n}^{c r} \\
\tau_{n}^{c r}
\end{array}\right], \quad t_{n}^{c r} \in \mathbb{C}^{\left(n-m_{c r}\right) \times m_{c r}}, \quad \tau_{n}^{c r} \in \mathbb{C}^{m_{c r} \times m_{c r}}
$$

In the other case we get

$$
L_{n} Y_{n}=\left[\begin{array}{l}
t_{n}^{c r} \\
\tilde{\tau}_{n}^{c r} \\
0
\end{array}\right]=\left[\begin{array}{l}
t_{n-1}^{c r} \\
\tau_{n-1}^{c r} \\
0
\end{array}\right]=\left[\begin{array}{c}
L_{n-1} Y_{n-1} \\
0
\end{array}\right] \quad \text { with } \quad \tilde{\tau}_{n}^{c r} \in \mathbb{C}^{\left(m_{c r}-1\right) \times m_{c r}}
$$

At a deflation step, the residual block $R_{n}^{c r}$ corresponding to the block iterate $X_{n}^{c r}$ in (60) is given by

$$
R_{n}^{c r}=V_{n-1}\left[\begin{array}{l}
0  \tag{64}\\
\tilde{\tau}_{n}^{c r}
\end{array}\right]
$$

The matrices $\hat{V}_{n}^{c r}$ and $\hat{V}_{n}^{d f}$ consist only of zero column vectors. Since $\tilde{\tau}_{n}^{c r}$ is a rectangular matrix, there exist a vector $\gamma$ with $\tilde{\tau}_{n}^{c r} \gamma=0$. Multiplying (64) from the right by $\gamma$ gives

$$
R_{n}^{c r} \gamma=V_{n-1}\left[\begin{array}{l}
0  \tag{65}\\
\tilde{\tau}_{n}^{c r}
\end{array}\right] \gamma=0 .
$$

In practice, the vector $\gamma$ is computed by means of the last $m_{c r}$ columns and the last $m_{c r}-1$ rows of the matrix $L_{n-1}$.
Using (65), we can express the approximate solution of the single linear system $A x_{\gamma}=B^{c r} \gamma$ in terms of a linear combination of columns of $X_{n}^{c r}$ by setting

$$
\begin{equation*}
x_{\gamma}=X_{n}^{c r} \gamma . \tag{66}
\end{equation*}
$$

Consequently, it is possible to delete one of the linear systems from block iterations, e.g., one can delete the $j$-th system. The solution vector corresponding to the $j$-th system can be constructed when all vectors in the updated block iterate $X_{n}^{c r}$ have converged after $n_{\text {max }}$ iterations. Using (66), we get

$$
\begin{equation*}
x_{n_{\text {max }}}^{\left(i_{j}\right)}=\frac{1}{\gamma_{j}}\left(x_{\gamma}-\sum_{\substack{k=1 \\ k \neq j}}^{m_{c r}} x_{n_{\max }}^{\left(i_{k}\right)} \gamma_{k}\right) . \tag{67}
\end{equation*}
$$

Now, we consider shifted systems of the type (41) with (42). Using

$$
V_{n}=P_{n}^{(j)} U_{n}^{(j)} \quad \text { for all } \quad 1 \leq j \leq n
$$

then the $n$-th iterate $x_{n}^{(j)}$ for the $j$-th system is defined by

$$
\begin{align*}
x_{n}^{(j)} & =x_{0}^{(j)}+V_{n} z_{n}^{(j)}, \quad z_{n}^{(j)} \in \mathbb{C}^{n}, \\
& =x_{0}^{(j)}+P_{n}^{(j)} U_{n}^{(j)} z_{n}^{(j)} \\
& =x_{0}^{(j)}+P_{n}^{(j)} y_{n}^{(j)}, \quad y_{n}^{(j)}=U_{n}^{(j)} z_{n}^{(j)} \in \mathbb{C}^{n}, \\
& =x_{n-1}^{(j)}+p_{n}^{(j)}\left(e_{n}^{T} y_{n}^{(j)}\right) . \tag{68}
\end{align*}
$$

Using (54) and $A^{(j)}=A+\sigma_{j} I$, we get

$$
\begin{equation*}
T_{n}^{(j)}=V_{n}^{T} A^{(j)} V_{n}=V_{n}^{T}\left(A+\sigma_{j} I\right) V_{n}=T_{n}+\sigma_{j} I . \tag{69}
\end{equation*}
$$

Using the decomposition $T_{n}^{(j)}=L_{n}^{(j)} U_{n}^{(j)}, y_{n}^{(j)}$ is the unique solution of the least-squares problem

$$
y_{n}^{(j)}=\arg \min _{y \in \mathbb{C}^{n}}\left\|\left[\begin{array}{l}
\Theta_{j}^{c r}  \tag{70}\\
0
\end{array}\right]-L_{n}^{(j)} y\right\| .
$$

### 5.3 Preconditioned System of Linear Algebraic Equations

Unfortunately, standard preconditioning techniques with a preconditioner

$$
\begin{equation*}
M=M_{1} M_{2}=\left(M_{1} M_{2}\right)^{T}=M^{T} \tag{71}
\end{equation*}
$$

such as SSOR preconditioning, destroy the special structure when they are applied to shifted linear systems. The only technique we are aware of that allows to preserve the shifted structure is polynomial preconditioning (see [7]).

Using standard preconditioning techniques, we apply the symmetric band Lanczos method to the shifted linear systems (41)

$$
\begin{equation*}
\tilde{A}^{(j)} \tilde{x}^{(j)}=\tilde{b}^{(j)}, \quad j=1, \ldots, n_{b}, \tag{72}
\end{equation*}
$$

with

$$
\begin{gather*}
\tilde{A}^{(j)}=M_{1}^{-1}\left(A+\sigma_{j} M_{1} M_{2}\right) M_{2}^{-1}=M_{1}^{-1} A M_{2}^{-1}+\sigma_{j} I, \\
\tilde{b}^{(j)}=M_{1}^{-1} b^{(j)} \quad \text { and } \quad \tilde{x}^{(j)}=M_{2} x^{(j)} . \tag{73}
\end{gather*}
$$

It is easy to see that the linear systems (41) and (72) with (73) are not equivalent. Only, for $M=I$ the systems (41) and (72) are equivalent.

It is common when implementing algorithms which involve a two-sided preconditioner to avoid the use of the $\tilde{x}^{(j)}$ variable. It is possible to write the resulting algorithm in terms of quantities corresponding to the system (41). We have the following analogies:

$$
\begin{array}{ll}
\hline v_{n} \rightarrow \tilde{v}_{n}=M_{1}^{-1} v_{n}, & p_{n} \rightarrow \tilde{p}_{n}=M_{2} p_{n}, \\
\tilde{v}_{n}^{T} \tilde{v}_{n}=v_{n}^{T} M^{-1} v_{n}, & \tilde{p}_{n}^{T} \tilde{p}_{n}=p_{n}^{T} M p_{n}, \\
\tilde{v}_{n}^{T} \tilde{A} \tilde{v}_{n}=v_{n}^{T} M^{-1} A M^{-1} v_{n}, & \tilde{p}_{n}^{T} \tilde{A} \tilde{p}_{n}=p_{n}^{T} A p_{n} . \\
\hline
\end{array}
$$

### 5.4 Algorithm

We summarize the basic structure of the symmetric band Lanczos QMR method with deflation, but without look-ahead (see [8, 4]).
0. Input: $A=A^{T}, B=\left[b^{(1)} b^{(2)} \ldots b^{\left(n_{b}\right)}\right],\left\{\sigma_{1}, \ldots, \sigma_{n_{b}}\right\}$.

A deflation tolerance dtol.
Choose $X_{0}$ and set $R=\left[r^{(1)} r^{(2)} \ldots r^{\left(n_{b}\right)}\right]=R_{0}=B-A X_{0}$.
For $j=1, \ldots, n_{b}$, set $\hat{v}_{j}=r^{(j)}$. Set $m_{c r}=n_{b}$ and index set $\mathcal{I}=\emptyset$.
For $n=1,2, \ldots$, do:

1. Compute $\beta=\left\|M_{1}^{-1} \hat{v}_{n}\right\|$. Set $\mu=n-m_{c r}$.

Decide if $\hat{v}_{n}$ should be deflated. If yes, i.e., $\beta<d t o l$, do the following:
1.1 If $\mu \leq 0$, then delete one column vector from $X_{0}^{c r}, R_{0}^{c r}$, and $\Theta^{c r}$.
1.2 If $\mu>0$ and the deflated vector $\hat{v}_{n}$ is nonzero, then set $\mathcal{I}=\mathcal{I} \cup\{\mu\}$. Delete one column vector from $X_{n}^{c r}$ and $R_{n}^{c r}$.
1.3 Set $m_{c r}=m_{c r}-1$. If $m_{c r}=0$, set $n=n-1$ and stop.

For $j=n, n+1, \ldots, n+m_{c r}-1$, set $\hat{v}_{j}=\hat{v}_{j+1}$.
Return to Step 1.
2. Update $V$ and preconditioned $V$ sequences.
2.1 Normalize $\hat{v}_{n}$ to become the $n$-th Lanczos vector: $v_{n}=\frac{\hat{v}_{n}}{\beta}$.
2.2 If $\mu \leq 0$, then set $\theta_{n, n}^{c r}=\beta$.
2.3 If $\mu>0$, then set $t_{n, \mu}=\beta$.
2.4 Compute $\delta_{n}=v_{n}^{T} M^{-1} v_{n}$. If $\delta_{n}=0$ stop.
2.5 Advance the right Krylov subspace: $\quad \hat{v}_{n+m_{c r}}=A\left(M^{-1} v_{n}\right)$.
2.6 Set $k=\max \{1, \mu\}$. Orthogonalize the vector $\hat{v}_{n+m_{c r}}$ against vectors $v_{j}$, $j \in \mathcal{I} \cup\{k, \ldots, n-1\}:$

$$
t_{j, n}=\frac{v_{j}^{T} M^{-1} \hat{v}_{n+m_{c r}}}{\delta_{j}} \quad \text { and } \quad \hat{v}_{n+m_{c r}}=\hat{v}_{n+m_{c r}}-v_{j} t_{j, n}
$$

2.7 Orthogonalize the vectors $\hat{v}_{n+j}, j=1, \ldots, m_{c r}$, against vector $v_{n}$ :

If $\mu+j \leq 0$, then

$$
\theta_{n, n+j}^{c r}=\frac{v_{n}^{T} M^{-1} \hat{v}_{n+j}}{\delta_{n}} \quad \text { and } \quad \hat{v}_{n+j}=\hat{v}_{n+j}-v_{n} \theta_{n, n+j}^{c r}
$$

If $\mu+j>0$, then

$$
t_{n, \mu+j}=\frac{v_{n}^{T} M^{-1} \hat{v}_{n+j}}{\delta_{n}} \quad \text { and } \quad \hat{v}_{n+j}=\hat{v}_{n+j}-v_{n} t_{n, \mu+j}
$$

For all $j=1, \ldots, n_{b}$ for which $x_{n}^{(j)}$ has not converged yet:
3. Compute the decomposition of the matrix $T_{n}^{(j)}$ :

$$
T_{n}^{(j)}=T_{n}+\sigma_{j} I=L_{n}^{(j)} U_{n}^{(j)}
$$

4. Compute the vector $p_{n}^{(j)}$ :

$$
P_{n}^{(j)}=\left[p_{1}^{(j)} p_{2}^{(j)} \ldots p_{n}^{(j)}\right]=V_{n}\left(U_{n}^{(j)}\right)^{-1}, \quad p_{n}^{(j)}=M^{-1} v_{n}-\sum_{i=\max \{1, \mu\}}^{n-1} p_{i}^{(j)} u_{i, n}^{(j)}
$$

5. Compute the solution $y_{n}^{(j)}$ of the least-squares problem (70).
6. Compute the $n$-th iterate $x_{n}^{(j)}$ (see (68)).

## End for $(j)$.

7. Check if all solution vectors $x_{n}^{(j)}$ have converged. If yes, then recover solution vectors (see (67)) and stop.

End for ( $n$ ).

## 6 Polynomial Preconditioning

We consider the shifted linear system

$$
\begin{equation*}
A^{(\sigma)} x=(A+\sigma I) x=b \tag{74}
\end{equation*}
$$

We use polynomial preconditioning to speed up the convergence of the iterative methods for the solution of (74), i.e.,

$$
\begin{equation*}
\mathcal{P}_{m}^{(\sigma)}\left(A^{(\sigma)}\right) A^{(\sigma)} x=\mathcal{P}_{m}^{(\sigma)}\left(A^{(\sigma)}\right) b \tag{75}
\end{equation*}
$$

for left preconditioning and

$$
\begin{equation*}
A^{(\sigma)} \mathcal{P}_{m}^{(\sigma)}\left(A^{(\sigma)}\right) y=\mathcal{P}_{m}^{(\sigma)}\left(A^{(\sigma)}\right) A^{(\sigma)} y=b, \quad x=\mathcal{P}_{m}^{(\sigma)}\left(A^{(\sigma)}\right) y \tag{76}
\end{equation*}
$$

for right preconditioning, respectively. Here, $\mathcal{P}_{m}^{(\sigma)}$ is a suitable chosen polynomial in $A^{(\sigma)}$ of a small degree, i.e., with degree no more than $m$. Both linear systems (75) and (76) are equivalent. We seek a polynomial $\mathcal{P}_{m}^{(\sigma)}$ with the following two properties [6, 5]:

- The coefficient matrix $\mathcal{P}_{m}^{(\sigma)}\left(A^{(\sigma)}\right) A^{(\sigma)}$ is again a shifted matrix.
- $\mathcal{P}_{m}^{(\sigma)}\left(A^{(\sigma)}\right)$ is an optimal polynomial preconditioner, i.e., the convergence of the iterative method, applied to the preconditioned system, is speed up optimally.

First, for any polynomial, we can represent $A^{(\sigma)} \mathcal{P}_{m}^{(\sigma)}\left(A^{(\sigma)}\right)$ in the form

$$
\begin{equation*}
A^{(\sigma)} \mathcal{P}_{m}^{(\sigma)}\left(A^{(\sigma)}\right)=(A+\sigma I) \mathcal{P}_{m}^{(\sigma)}(A+\sigma I)=A \mathcal{P}_{m}(A)+\tau I \tag{77}
\end{equation*}
$$

with $\tau \in \mathbb{C}$. Note that $\mathcal{P}_{m}^{(\sigma)}, \mathcal{P}_{m}$, and $\tau$ are related by

$$
\begin{equation*}
(z+\sigma) \mathcal{P}_{m}^{(\sigma)}(z+\sigma)=z \mathcal{P}_{m}(z)+\tau \quad \text { and } \quad \tau=\sigma \mathcal{P}_{m}(-\sigma) \tag{78}
\end{equation*}
$$

We note that the coefficient matrix $A \mathcal{P}_{m}(A)$ of the preconditioned system (77) is Hermitian if, and only if, $\mathcal{P}_{m}$ is a real polynomial. In order to guarantee that $A \mathcal{P}_{m}(A)$ is nonsingular, we require that $\mathcal{P}_{m}(z) \neq 0$ for all $z \in S$ with

$$
\varrho(A) \subseteq S=[a, b] \cup[c, d], \quad c<d<0<a<b
$$

where $\varrho(A)$ is the spectrum of $A$.
Next, we turn to the question of optimal choice of polynomial $\mathcal{P}_{m}$. We have two different cases:

- $z \mathcal{P}_{m}(z)>0 \quad \forall z \in S$
- $z \mathcal{P}_{m}(z)>0 \quad \forall z \in[a, b] \quad$ and $\quad z \mathcal{P}_{m}(z)<0 \quad \forall z \in[c, d]$.

If the last case holds, then the preconditioned system remains indefinite. We can now state the main result in the following form [6]:

Let $S=[a, b] \cup[c, d]$ be the union of a positive and negative interval with $c<$ $d<0<a<b$ and $\Gamma=\{(\gamma, \delta) \in \mathbb{R} \times \mathbb{R}: \delta>0\}$ a parameter set. The optimal polynomial $\mathcal{P}_{m}^{*}(z)$ of

$$
\begin{gather*}
w(\gamma, \delta)=\min _{\mathcal{P}_{m}}\left\|f-z \mathcal{P}_{m}\right\|_{g} \\
\left\|f-z \mathcal{P}_{m}\right\|_{g}=\max _{z \in S}\left|g(z)\left(f(z)-z \mathcal{P}_{m}(z)\right)\right|, \tag{79}
\end{gather*}
$$

where

$$
g(z)=\left\{\begin{array}{lll}
1 & \text { if } & z>0 \\
\delta & \text { if } & z<0
\end{array}, \quad f(z)=\left\{\begin{array}{lll}
1 & \text { if } & z>0 \\
\gamma & \text { if } & z<0
\end{array}\right.\right.
$$

is an indefinite polynomial preconditioner with

$$
\gamma=\frac{\bar{d}+\bar{c}}{\bar{b}+\bar{a}} \quad \text { and } \quad \delta=\frac{\bar{b}-\bar{a}}{\bar{d}-\bar{c}}
$$

The numbers $\bar{a}, \bar{b}, \bar{c}$, and $\bar{d}$ are defined by

$$
\begin{array}{ll}
\bar{a}=\min _{z \in[a, b]} z \mathcal{P}_{m}(z), & \bar{b}=\max _{z \in[a, b]} z \mathcal{P}_{m}(z), \\
\bar{c}=\min _{z \in[c, d]} z \mathcal{P}_{m}(z), & \bar{d}=\max _{z \in[c, d]} z \mathcal{P}_{m}(z) .
\end{array}
$$

Moreover, there exist parameters $\gamma_{0}$ and $\delta_{0},\left(\gamma_{0}, \delta_{0}\right) \in \Gamma$, such that $\mathcal{P}_{m}^{*}\left(z, \gamma_{0}, \delta_{0}\right)$ is an optimal indefinite polynomial preconditioner.
(79) is a linear Chebyshev approximation problem depending on the two parameters $(\gamma, \delta) \in \Gamma$. We seek to approximate $f(z)$ by polynomials of the form $z \mathcal{P}_{m}(z)$ in the weighted uniform norm $\|\cdot\|_{g}$. The standard tool for the numerical solution of such general real Chebyshev approximation problems is the method of Remez. The Remez type procedure is based on the equioscillation property [6]. Implementation details are described in [19].

## 7 Initial Guesses

The use of initial guesses for systems with multiple shifts is not as simple. The problem of initial guesses is related to the more general problem of solving systems with multiple shifts each with a different right-hand side (see [17]).

### 7.1 Multiple Shifts with Identical Right-Hand Sides

We consider the systems (41) with the same right-hand sides (see (28)). Using the initial guesses, $x_{0}^{(j)}$, we get the residuals

$$
\begin{equation*}
r_{0}^{(j)}=b-\left(A+\sigma_{j} I\right) x_{0}^{(j)}, \quad j=1, \ldots, n_{b} \tag{80}
\end{equation*}
$$

and then solve

$$
\begin{align*}
b-\left(A+\sigma_{j} I\right) x_{0}^{(j)} & =r_{0}^{(j)} \\
\left(A+\sigma_{j} I\right) x^{(j)}-\left(A+\sigma_{j} I\right) x_{0}^{(j)} & =r_{0}^{(j)} \\
\left(A+\sigma_{j} I\right)\left(x^{(j)}-x_{0}^{(j)}\right) & =r_{0}^{(j)} \\
\left(A+\sigma_{j} I\right) z^{(j)} & =r_{0}^{(j)}, \quad j=1, \ldots, n_{b} . \tag{81}
\end{align*}
$$

Then, the solutions are given by $x^{(j)}=x_{0}^{(j)}+z^{(j)}$ for all $j=1, \ldots, n_{b}$. In general, the new right-hand sides, $r_{0}^{(j)}$, are not collinear. Therefore, the Krylov subspaces, $\mathcal{K}_{n}\left(A+\sigma_{j} I, r_{0}^{(j)}\right)$, are not the same. Take

$$
\begin{equation*}
x_{0}^{(j)}=\left\{\prod_{\substack{i=1 \\ i \neq j}}^{n_{b}}\left(A+\sigma_{i} I\right)\right\} w \tag{82}
\end{equation*}
$$

for any vector $w$, then the new right-hand sides are all equal to

$$
\begin{equation*}
r_{0}^{(j)}=b-\left(A+\sigma_{j} I\right)\left\{\prod_{\substack{i=1 \\ i \neq j}}^{n_{b}}\left(A+\sigma_{i} I\right)\right\} w=b-\left\{\prod_{i=1}^{n_{b}}\left(A+\sigma_{i} I\right)\right\} w . \tag{83}
\end{equation*}
$$

The $n_{b}$ systems (81) with (82) and (83) can be solved with standard multi-shift Krylov methods [19]. The problem is to find the best vector $w$ for (82).

Given approximate solutions

$$
\begin{equation*}
\tilde{x}^{(j)} \approx\left(A+\sigma_{j} I\right)^{-1} b, \quad j=1, \ldots, n_{b} . \tag{84}
\end{equation*}
$$

Using (83), we get

$$
b \approx\left\{\prod_{i=1}^{n_{b}}\left(A+\sigma_{i} I\right)\right\} w
$$

Next, from the above expression, we derive the vector $w$ :

$$
\begin{aligned}
\tilde{x}^{(j)} \approx\left(A+\sigma_{j} I\right)^{-1} b & \approx\left\{\prod_{\substack{i=1 \\
i \neq j}}^{n_{b}}\left(A+\sigma_{i} I\right)\right\} w \\
& \approx\left(A+\sigma_{k} I\right)\left\{\prod_{\substack{i=1 \\
i \neq j, k}}^{n_{b}}\left(A+\sigma_{i} I\right)\right\} w, \quad 1 \leq j \leq n_{b},
\end{aligned}
$$

$$
\begin{aligned}
& \tilde{x}^{(k)} \approx\left(A+\sigma_{k} I\right)^{-1} b \approx\left\{\prod_{\substack{i=1 \\
i \neq k}}^{n_{b}}\left(A+\sigma_{i} I\right)\right\} w \\
& \approx\left(A+\sigma_{j} I\right)\left\{\prod_{\substack{i=1 \\
i \neq j, k}}^{n_{b}}\left(A+\sigma_{i} I\right)\right\} w, \quad 1 \leq k \leq n_{b}, \\
& \tilde{x}^{(l)} \approx\left(A+\sigma_{l} I\right)^{-1} b \approx\left(A+\sigma_{j} I\right)\left\{\prod_{\substack{i=1 \\
i \neq j, l}}^{n_{b}}\left(A+\sigma_{i} I\right)\right\} w, \quad 1 \leq l \leq n_{b}, \\
& \Downarrow \\
& \frac{\tilde{x}^{(j)}-\tilde{x}^{(k)}}{\sigma_{k}-\sigma_{j}} \approx\left\{\prod_{\substack{i=1 \\
i \neq j, k}}^{n_{b}}\left(A+\sigma_{i} I\right)\right\} w \\
& \approx\left(A+\sigma_{l} I\right)\left\{\prod_{\substack{i=1 \\
i \neq j, k, l}}^{n_{b}}\left(A+\sigma_{i} I\right)\right\} w, \quad 1 \leq j, k \leq n_{b}, \\
& \frac{\tilde{x}^{(j)}-\tilde{x}^{(l)}}{\sigma_{l}-\sigma_{j}} \approx\left\{\prod_{\substack{i=1 \\
i \neq j, l}}^{n_{b}}\left(A+\sigma_{i} I\right)\right\} w \\
& \approx\left(A+\sigma_{k} I\right)\left\{\prod_{\substack{i=1 \\
i \neq j, k, l}}^{n_{b}}\left(A+\sigma_{i} I\right)\right\} w, \quad 1 \leq j, l \leq n_{b}, \\
& \Downarrow \\
& \frac{\tilde{x}^{(j)}}{\left(\sigma_{k}-\sigma_{j}\right)\left(\sigma_{l}-\sigma_{j}\right)}+ \\
& \frac{\tilde{x}^{(k)}}{\left(\sigma_{j}-\sigma_{k}\right)\left(\sigma_{l}-\sigma_{k}\right)}+ \\
& \frac{\tilde{x}^{(l)}}{\left(\sigma_{j}-\sigma_{l}\right)\left(\sigma_{k}-\sigma_{l}\right)} \approx\left\{\prod_{\substack{i=1 \\
i \neq j, k, l}}^{n_{b}}\left(A+\sigma_{i} I\right)\right\} w, \quad 1 \leq j, k, l \leq n_{b} .
\end{aligned}
$$

This process can be continued up to $\prod_{i}^{n_{b}}\left(A+\sigma_{i} I\right) \equiv 1$. Using (84), we get

$$
\begin{equation*}
w=\sum_{j=1}^{n_{b}} c_{j} \tilde{x}^{(j)} \quad \text { with } \quad c_{j}=\prod_{\substack{i=1 \\ i \neq j}}^{n_{b}} \frac{1}{\sigma_{i}-\sigma_{j}} \tag{85}
\end{equation*}
$$

Note that the coefficients $c_{j}$ can be become very large if the shifts $\sigma_{j}$ with smaller differences increase.

### 7.2 Multiple Shifts with Multiple Sources

Next, we are interested in solving the systems (41) with multiple sources. We need to find a set of initial values, $x_{0}^{(j)}, j=1, \ldots, n_{b}$, that give a common residul $r_{0}$ :

$$
r_{0}=r_{0}^{(j)}=b^{(j)}-\left(A+\sigma_{j} I\right) x_{0}^{(j)}, \quad j=1, \ldots, n_{b}
$$

This can be solved by setting

$$
x_{0}^{(j)}=\sum_{k=0}^{n_{b}-2} A^{k} z^{(j, k)} .
$$

Then we get

$$
\begin{aligned}
r_{0} & =b^{(j)}-\left(A+\sigma_{j} I\right) x_{0}^{(j)} \\
& =b^{(j)}-\left(A+\sigma_{j} I\right)\left\{\sum_{k=0}^{n_{b}-2} A^{k} z^{(j, k)}\right\} \\
& =b^{(j)}-\left\{\sum_{k=0}^{n_{b}-2} A^{k+1} z^{(j, k)}+\sigma_{j} \sum_{k=0}^{n_{b}-2} A^{k} z^{(j, k)}\right\}, j=1, \ldots, n_{b} .
\end{aligned}
$$

Equating powers of $A$ and solve for $z^{(j, k)}$ in terms of $b^{\prime} s$ gives the following equations for $i, j=$ $1, \ldots, n_{b}$ :

$$
\begin{aligned}
z^{\left(j, n_{b}-2\right)} & =z^{\left(i, n_{b}-2\right)} \\
z^{(j, k-1)}+\sigma_{j} z^{(j, k)} & =z^{(i, k-1)}+\sigma_{i} z^{(i, k)}, \quad k=n_{b}-2, \ldots, 1 \\
b^{(j)}-\sigma_{j} z^{(j, 0)} & =b^{(i)}-\sigma_{i} z^{(i, 0)}
\end{aligned}
$$

One can also solve this by considering the polynomials

$$
q_{j}(A)=\left(A+\sigma_{j} I\right) p_{j}(A) \quad \text { with } \quad x_{0}^{(j)}=\sum_{k=0}^{n_{b}-2} A^{k} z^{(j, k)} \equiv p_{j}(A)
$$

at the special cases $A=-\sigma_{l} I, 1 \leq l \leq n_{b}$ :

$$
r_{0}^{\sigma_{l}}=b^{(l)}-q_{l}\left(-\sigma_{l} I\right)=b^{(l)}-\left(-\sigma_{l} I+\sigma_{l} I\right) x_{0}^{(l)}=b^{(l)} .
$$

This gives the $n_{b}$ equations for fixed $j$

$$
q_{j}\left(-\sigma_{l} I\right)=b^{(j)}-r_{0}^{\sigma_{l}}=b^{(j)}-b^{(l)}, 1 \leq l \leq n_{b} .
$$

Since $q_{j}(A)$ is a polynomial of order $n_{b}-1$ in $A$ the system is uniquely determined. The polynomial

$$
\begin{equation*}
q_{j}(A)=\sum_{k=1}^{n_{b}}\left\{\prod_{\substack{i=1 \\ i \neq k}}^{n_{b}} \frac{A+\sigma_{i} I}{\sigma_{i}-\sigma_{k}}\right\}\left(b^{(j)}-b^{(k)}\right) \tag{86}
\end{equation*}
$$

satisfies these equations. Using $x_{0}^{(j)}=p_{j}(A)=\left(A+\sigma_{j} I\right)^{-1} q_{j}(A)$ we get

$$
\begin{equation*}
x_{0}^{(j)}=\sum_{\substack{k=1 \\ k \neq j}}^{n_{b}}\left\{\prod_{\substack{i=1 \\ i \neq j, k}}^{n_{b}} \frac{A+\sigma_{i} I}{\sigma_{i}-\sigma_{k}}\right\}\left(\frac{b^{(j)}-b^{(k)}}{\sigma_{j}-\sigma_{k}}\right) . \tag{87}
\end{equation*}
$$

## 8 Numerical Results

A nonequidistant mesh of 57664 elementary cells including graded PML regions is used for the discretization of (4), that means the order of the system of linear algebraic equations is 172992 (see (19)). The number of internal points $y_{I}^{(G)}$ (see Table 1 and Eq. (20)) is 152608 . We use three different matrices $D$ (see (23)):

$$
\text { (I) } D=\operatorname{diag}(A), \quad \text { (II) } \quad D=\frac{1}{10} \cdot \operatorname{diag}(A), \quad \text { (III) } \quad D=10 \cdot \operatorname{diag}(A) \text {. }
$$

We analyze the following systems of linear equations:

| Linear system |  |  | I | II | III |
| :--- | :--- | :--- | :--- | :--- | :--- |
| 0 | $A x^{(j)}=b^{(j)}$ | $\left(\alpha_{j}=1, \beta_{j}=0\right)$ | $I_{0}$ |  |  |
| 1 | $\left(\alpha_{j} A+\beta_{j} D\right) x^{(j)}=b^{(j)}$ | (see Sec. 5.2) | $I_{1}$ | $I I_{1}$ | $I I I_{1}$ |
| 2 | $\left(\alpha_{j} A+\beta_{j} D\right) x^{(j)}=b^{(j)}$ | (see Sec. 5.2) | $I_{2}$ | $I I_{2}$ | $I I I_{2}$ |
| 3 | $\left(\alpha_{j} A+\beta_{j} D\right) x^{(j)}=b^{(j)}$ | (see Sec. 7.2) | $I_{3}$ | $I I_{3}$ | $I I I_{3}$ |

The systems of linear equations in row $2\left(I_{2}, I I_{2}, I I I_{2}\right)$ are computed separately. The systems of linear equations in row $3\left(I_{3}, I I_{3}, I I I_{3}\right)$ are computed with new initial guesses (see (87)) and a common new right-hand side (see Sec. 7.2). This method does not work very good for more and/or smaller shifts. The initial residuals may be large. Thus, we use a fixed number of right-hand sides, $n b m$, to get new initial guesses and a new right-hand side:

$$
I_{3}: \quad n b m=6, \quad I I_{3}: \quad n b m=3, \quad I I I_{3}: \quad n b m=10
$$

The right-hand sides $b^{(j)}, j=1, \ldots, n_{b}$, are generated by an uniformly distributed random number generator. The stopping criterion was a reduction of the norm of the residual for the preconditioned system (41-43) by $10^{-8}$.

Based on the family of approximation problems (79), we have computed indefinite polynomial preconditioners. For this purpose, the Remez algorithm in [19, Sec. 6.1] was used. Using the

Nelder-Mead method (see [19, Sec. 6.2]), optimal indefinite polynomial preconditioners were computed by solving the constrained optimization problem numerically. The experiments were conducted on a workstation with Intel Xeon CPU W3520 ( $2,67 \mathrm{GHz}, 8$ cores) running 64-bit Linux system.

We compare the iteration counts required to solve the $n_{b} \in\{1,3,5,10\}$ linear systems (see (29) and (41)) using the symmetric band Lanczos method without polynomial preconditioning with the number of iterations required to solve with polynomial preconditioning. The choice of the parameters $(\gamma, \delta) \in \Gamma$ and the computation of the 'optimal' parameters $\left(\gamma_{0}, \delta_{0}\right)$ is described in [19].
In Table 2 are denoted the estimated boundaries for the sets $S$, the parameters $(\gamma, \delta) \in \Gamma$ for the examples considered here, and also the computed 'optimal' parameters $\left(\gamma_{0}, \delta_{0}\right) \in \Gamma$. The degree of the Lagrange polynomial $L_{j}(z)$ is $l=8$ and $l=9$, respectively [19, Eq. (51)].

Table 2: The boundaries for $S$

I

| $a$ | $b$ | $c$ | $d$ |
| :---: | :---: | :---: | :---: |
| $2.24 \cdot 10^{-6}$ | 2.00 | -0.45 | $-5.20 \cdot 10^{-7}$ |
| $\gamma \approx-0.2241, \quad \delta \approx 4.4610$ |  |  |  |

$l=9: \quad \gamma_{0} \approx-0.7339, \quad \delta_{0} \approx 4.3407$

II
$l=9: \quad \gamma_{0} \approx-0.7339, \quad \delta_{0} \approx 4.3407$

| a | b | c | d |
| :---: | :---: | :---: | :---: |
| $1.27 \cdot 10^{-4}$ | 20.00 | -3.66 | $-5.19 \cdot 10^{-6}$ |
| $\gamma \approx-0.1825, \quad \delta \approx 5.4782$ |  |  |  |

$l=8: \quad \gamma_{0} \approx-0.5635, \quad \delta_{0} \approx 5.5894$

III

| a | b | c | d |
| :---: | :---: | :---: | :---: |
| $9.48 \cdot 10^{-7}$ | 0.20 | $-6.48 \cdot 10^{-2}$ | $-5.19 \cdot 10^{-8}$ |
| $\gamma \approx-0.3240, \quad \delta \approx 3.0864$ |  |  |  |
| $\quad l=9:$ | $\gamma_{0} \approx-0.5066, \quad \delta_{0} \approx 3.2273$ |  |  |

We choose for the preconditioner $\mathcal{P}_{m}(A)$ (see (77)) the linear case. This choice decreases the numerical effort and is more stable. The Tables $3-6$ shows the numbers of iterations for shifted and unshifted linear system, respectively. A convergence history of the three examples (see Tables 4-6) can be found graphically in the Figs. 4-6.

Table 3: Number of iterations for unshifted matrices, $I_{0}$

| Number of unshifted linear systems, $n_{b}$ | 1 | 3 | 5 | 10 |
| ---: | :--- | ---: | ---: | ---: | ---: |
| No preconditioning | 1128 | 2712 | 4355 | 8844 |
| Preconditioning: $\left(\gamma_{0}, \delta_{0}\right)$ | 569 | 1219 | 1965 | 4385 |

Table 4: Number of iterations for shifted matrices, $I_{1}-I I I_{1}$

| $I_{1}$ | Number of shifted linear systems, $n_{b}$ | 1 | 3 | 5 | 10 |
| :---: | :---: | :---: | :---: | :---: | :---: |
|  | No preconditioning | 119 | 10722 | 18064 | 36212 |
|  | Preconditioning: $\left(\gamma_{0}, \delta_{0}\right)$ | 58 | 5374 | 9225 | 18784 |
| $I I_{1}$ | Number of shifted linear systems, $n_{b}$ | 1 | 3 | 5 | 10 |
|  | No preconditioning | 369 | 13329 | 14209 | 76780 |
|  | Preconditioning: $\left(\gamma_{0}, \delta_{0}\right)$ | 195 | 4726 | 5657 | 34868 |
| $I I I_{1}$ | Number of shifted linear systems, $n_{b}$ | 1 | 3 | 5 | 10 |
|  | No preconditioning | 41 | 3642 | 6103 | 12581 |
|  | Preconditioning: $\left(\gamma_{0}, \delta_{0}\right)$ | 20 | 1716 | 2953 | 6084 |

## 9 Conclusions

We have analyzed the solution of linear systems of equations with multiple right-hand sides each with a different shift. We have presented a variant of the symmetric band Lanczos process with multiple starting vectors. The symmetric band Lanczos process is based on coupled recurrences. Polynomial preconditioners for indefinite linear systems leads to indefinite preconditioned coefficient matrices. Such polynomials can be obtained via the solution of linear Chebyshev approximation problems.
We have presented a method that produce initial guesses with a common right-hand side. It could be very useful. The main difficulty then is keeping the initial residual under control.
Another problem is to find suitable informations on the location of the eigenvalues of $A$, i.e., the bounds of the two intervals $[a, b]$ and $[c, d]$.

Table 5: Number of iterations for shifted matrices, $\mathrm{I}_{2}-\mathrm{III}_{2}$

| $I_{2}$ | Number of shifted linear systems, $n_{b}$ | 1 | 3 | 5 | 10 |
| :---: | :---: | :---: | :---: | :---: | :---: |
|  | No preconditioning | 119 | 4883 | 5683 | 13975 |
|  | Preconditioning: $\left(\gamma_{0}, \delta_{0}\right)$ | 58 | 2433 | 2823 | 6975 |
| $\mathrm{II}_{2}$ | Number of shifted linear systems, $n_{b}$ | 1 | 3 | 5 | 10 |
|  | No preconditioning | 369 | 6549 | 8826 | 33506 |
|  | Preconditioning: $\left(\gamma_{0}, \delta_{0}\right)$ | 195 | 3364 | 4569 | 18246 |
| $\mathrm{III}_{2}$ | Number of shifted linear systems, $n_{b}$ | 1 | 3 | 5 | 10 |
|  | No preconditioning | 41 | 2480 | 2746 | 3921 |
|  | Preconditioning: $\left(\gamma_{0}, \delta_{0}\right)$ | 20 | 1221 | 1343 | 2081 |

Table 6: Number of iterations for shifted matrices, $\mathrm{I}_{3}-\mathrm{III}_{3}$

| $I_{3}$ | Number of shifted linear systems, $n_{b}$ | 1 | 3 | 5 | 10 |
| :---: | :---: | :---: | :---: | :---: | :---: |
|  | No preconditioning | 115 | 3381 | 2970 | 5360 |
|  | Preconditioning: $\left(\gamma_{0}, \delta_{0}\right)$ | 56 | 1647 | 1575 | 2720 |
| $\mathrm{II}_{3}$ | Number of shifted linear systems, $n_{b}$ | 1 | 3 | 5 | 10 |
|  | No preconditioning | 364 | 4803 | 6551 | 22341 |
|  | Preconditioning: $\left(\gamma_{0}, \delta_{0}\right)$ | 178 | 2425 | 3318 | 11036 |
| $\mathrm{III}_{3}$ | Number of shifted linear systems, $n_{b}$ | 1 | 3 | 5 | 10 |
|  | No preconditioning | 37 | 1100 | 1135 | 1175 |
|  | Preconditioning: $\left(\gamma_{0}, \delta_{0}\right)$ | 20 | 581 | 594 | 637 |

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Figure 4: Number of iterations for (I)


Figure 5: Number of iterations for (II)
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Figure 6: Number of iterations for (III)

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