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# ALTERNATING ITERATIVE SCHEME FOR THE SOLUTION OF BLOCK-STRUCTURED SYSTEMS * 

Jan Šípek, Jan Zítko


#### Abstract

We consider the solution of linear system with a block-structured matrix of saddle point type. The solution technique is based on the idea of the classical alternatingdirection implicit iterative method where symmetric-antisymmetric splitting of the coefficient matrix is used. To find an optimal parameter for solving the system with a symmetric matrix, the polynomial filters are considered. The CGW method is used for systems with skew-symmetric matrix. The numerical tests compare the results obtained by using alternating iteration and GMRES and point out advantages of alternative iterations for larger systems.


## 1. Introduction

We consider solution of linear system of equations with the following block structure

$$
\underbrace{\left[\begin{array}{cc}
A & B^{T}  \tag{1}\\
-B & C
\end{array}\right]}_{\mathcal{A}} \underbrace{\left[\begin{array}{c}
u \\
p
\end{array}\right]}_{x}=\underbrace{\left[\begin{array}{c}
f \\
-g
\end{array}\right]}_{b},
$$

where $A \in \mathbb{R}^{n \times n}, B \in \mathbb{R}^{m \times n}, C \in \mathbb{R}^{m \times m}, f \in \mathbb{R}^{n}, g \in \mathbb{R}^{m}$, and $m<n$. Moreover, we assume that $\operatorname{rank}(B)=m$ and $N(A) \cap N(B)=\{0\}$. Systems of the form (1) arise in a number of scientific and engineering applications. There are obtained by calculation of the minimum of a functional with $m$ linear constrains or by the discretization of Navier-Stokes equations using the finite element or finite volume method; see References [1, 2, 14, 15]. In some special cases the system (1) has more interesting properties such as $A$ is symmetric positive definite (SPD) or $C=0$.

In our further investigation, we will assume that the matrix $A$ is non-symmetric with a positive definite symmetric part $\frac{1}{2}\left(A+A^{T}\right)$. In such a case we refer matrix $A$ as positive real. Define

$$
\begin{equation*}
\mathcal{G}:=\frac{1}{2}\left(\mathcal{A}+\mathcal{A}^{T}\right), \quad \mathcal{S}:=\frac{1}{2}\left(\mathcal{A}-\mathcal{A}^{T}\right) \tag{2}
\end{equation*}
$$

and similarly

$$
\begin{equation*}
G=\frac{1}{2}\left(A+A^{T}\right), \quad S=\frac{1}{2}\left(A-A^{T}\right) . \tag{3}
\end{equation*}
$$

[^0]It was proved (see [7]) that if $G$ is positive semidefinite, $\operatorname{rank}(B)=m$ and $N(G) \cap$ $N(B)=\{0\}$, then $\mathcal{A}$ is nonsingular, semipositive real and positive semistable (i.e. the eigenvalues of $\mathcal{A}$ have nonnegative real part). A number of methods have been proposed for the solution of (1) in the literature. Let us mention [12], [7], [2], [16], [10], [3], [18], [19], [20]. Very interesting technique based on an idea of the classic alternating direction implicit method described in Reference [13] appears in the last years.

Let $\alpha \in \mathbb{R}$ be a positive number. Write

$$
\mathcal{A}=(\mathcal{G}+\alpha \mathcal{I})+(\mathcal{S}-\alpha \mathcal{I})=(\mathcal{G}-\alpha \mathcal{I})+(\mathcal{S}+\alpha \mathcal{I})
$$

where $\mathcal{I}$ is identity matrix in $\mathbb{R}^{(n+m) \times(n+m)}$. Let us mention that

$$
\mathcal{G}+\alpha \mathcal{I}=\left[\begin{array}{cc}
G+\alpha I_{n} & 0  \tag{4}\\
0 & C+\alpha I_{m}
\end{array}\right], \quad \mathcal{S}+\alpha \mathcal{I}=\left[\begin{array}{cc}
S+\alpha I_{n} & B^{T} \\
-B & \alpha I_{m}
\end{array}\right] .
$$

According to the technique formulated for example in [13, 7], we write the following algorithm:

Algorithm 1: ADI method
Input: $x_{0}=\left[u^{0}, p^{0}\right]^{T}$ : an initial approximation.
Output: The solution $x$ of (1)
For $k=0$ till convergence do

$$
\begin{align*}
& (\mathcal{G}+\alpha \mathcal{I}) x^{k+\frac{1}{2}}=(\alpha \mathcal{I}-\mathcal{S}) x^{k}+b  \tag{G}\\
& (\mathcal{S}+\alpha \mathcal{I}) x^{k+1}=(\alpha \mathcal{I}-\mathcal{G}) x^{k+\frac{1}{2}}+b \tag{S}
\end{align*}
$$

end do

The matrix-vector multiplication on the right-hand side of $(G)$ and $(S)$ is evaluated only once for $k=0$. The right-hand side of each following "half-iteration" yields the left-hand side of previously executed one, after an easy manipulation.

Putting

$$
f_{1}^{k}=\alpha u^{k}-S u^{k}+f-B^{T} p^{k}, \quad g_{1}^{k}=\alpha p^{k}-g+B u^{k},
$$

we can see, that the first part of our algorithm involve the solution of two linear systems of the form

$$
\begin{equation*}
\left(G+\alpha I_{n}\right) u^{k+\frac{1}{2}}=f_{1}^{k}, \quad\left(C+\alpha I_{m}\right) p^{k+\frac{1}{2}}=g_{1}^{k} \tag{5}
\end{equation*}
$$

with positive definite matrices. Let us remark that the problems with $C=0$ are very often solved and in this case only one equation with positive definite matrix $G+\alpha I_{n}$ is solved. When solving the linear system (5) with conjugate gradient method, the smallest eigenvalues slow down the convergence. The same phenomenon has been observed by solving nonsymmetric systems using GMRES method and preconditioning
technique removing the small eigenvalues from the spectrum was studied by many authors. The situation in symmetric positive definite case is much more easier because the matrix is diagonalizable with real positive eigenvalues. This is the reason for splitting the coefficient matrix $\mathcal{A}$ into symmetric and skew-symmetric part and solve more systems of linear algebraic equations. For removing small eigenvalues, an invariant subspace associated with these ones is usually constructed. We have proposed the method of polynomial filters modified for symmetric systems. The authors [Giraud, Ruiz, Touhami, Arioli] proposed an algorithm which combine Chebyshev iteration with a block Lanczos procedure to accurately compute an orthonormal basis for the invariant subspace associated with the small eigenvalues of the matrix $A$. It uses Chebyshev polynomials to damp eigencomponents associated with the largest eigenvalues in the spectral decomposition of the considered vector. Both approaches for an efficient solving of symmetric system have the common idea to enforce conjugate gradient to work in the orthogonal complement of some invariant subspace associated with the smallest eigenvalues.

The second system in (S), rewritten according (4) as

$$
\begin{align*}
l c l\left(\alpha I_{n}+S\right) u^{k+1}+B^{T} p^{k+1} & =f_{2}^{k}  \tag{6}\\
-B u^{k+1}+\alpha p^{k+1} & =g_{2}^{k} \tag{7}
\end{align*}
$$

where

$$
f_{2}^{k}=\left(\alpha I_{n}-G\right) u^{k+\frac{1}{2}}+f, \quad g_{2}^{k}=\left(\alpha I_{m}-C\right) p^{k+\frac{1}{2}}-g,
$$

can be solved in various ways. The CGW method is often recommended for the solution of (6); see Reference [22]. An alternative approach [7] is to eliminate $u^{k+1}$ from the second equation using the first, which is useful for symmetric matrix $\mathcal{A}$ or to eliminate $p^{k+1}$ from the first equation using the second, which seems to be more useful for non-symmetric case. The first reduction described, leads on solution of following equation

$$
\begin{equation*}
\left[B\left(I_{n}+\alpha S\right)^{-1} B^{T}+\alpha I_{m}\right] p^{k+1}=B\left(I_{n}+\alpha S\right)^{-1} f_{2}^{k}+g_{2}^{k} . \tag{8}
\end{equation*}
$$

From (8) we get the solution vector $p^{k+1}$. Vector $u^{k+1}$ is defined by

$$
u^{k+1}=\left(\alpha I_{n}+S\right)^{-1}\left(f_{2}^{k}-B^{T} p^{k+1}\right)
$$

If $S=0$, then the system (8) degenerates to

$$
\begin{equation*}
\left(B B^{T}+\alpha^{2} I_{m}\right) p^{k+1}=B f_{2}^{k}+\alpha g_{2}^{k} \tag{9}
\end{equation*}
$$

and $u^{k+1}=\frac{1}{\alpha}\left(f_{2}^{k}-B^{T} p^{k+1}\right)$. Relation (8) is far to complicated in non-symmetric case and numerically too expensive. Hence, let us try to express $p^{k+1}$ :

$$
\begin{equation*}
p^{k+1}=\frac{1}{\alpha}\left(B u^{k+1}+g_{2}^{k}\right) \tag{10}
\end{equation*}
$$

and substitute to the first equation, from which we get $u_{k+1}$. Therefore we have

$$
\begin{equation*}
\left(\left(\alpha I_{n}+S\right)+\frac{1}{\alpha} B^{T} B\right) u^{k+1}=f_{2}^{k}-\frac{1}{\alpha} B^{T} g_{2}^{k} \tag{11}
\end{equation*}
$$

The aim of this paper is to compare several of these techniques in terms of numerical efficiency and to show the advantages of the alternating iterations.

The outline of this paper is as follows. In Section 2 the convergence and estimation of an optimal shift parameter is discussed. We devote Section 3 and 4 to application of polynomial filters for symmetric case. Various filtering techniques are shown for the construction of an invariant subspace. In Section 5 the CGW method is studied. In Section 6 the numerical results are shown. The comparison of alternating-iteration method with Krylov subspace methods is presented.

## 2. Convergence and estimation of optimal parameter $\alpha$

We are looking for such $\alpha$, that the method introduced, converge optimally. From $(\mathrm{G})$ and $(\mathrm{S})$ we become the following one step iterative method of the form

$$
x_{k+1}=(\mathcal{S}+\alpha \mathcal{I})^{-1}(\mathcal{G}-\alpha \mathcal{I})(\mathcal{G}+\alpha \mathcal{I})^{-1}(\mathcal{S}-\alpha \mathcal{I}) x_{k}+c(\alpha)
$$

Using the relation $\mathcal{A}=\mathcal{G}+\mathcal{S}$ we define following

$$
\begin{aligned}
\mathcal{M} & =(2 \alpha)^{-1}(\alpha \mathcal{I}+\mathcal{G})(\alpha \mathcal{I}+\mathcal{S}) \\
\mathcal{N} & =(2 \alpha)^{-1}(\alpha \mathcal{I}-\mathcal{G})(\alpha \mathcal{I}-\mathcal{S}) \\
\mathcal{W} & =\mathcal{I}-\mathcal{A} \mathcal{M}^{-1}=\mathcal{N} \mathcal{M}^{-1}=(\alpha \mathcal{I}-\mathcal{G})(\alpha \mathcal{I}-\mathcal{S})(\alpha \mathcal{I}+\mathcal{S})^{-1}(\alpha \mathcal{I}+\mathcal{G})^{-1}
\end{aligned}
$$

For the iteration $\left\{x_{k}\right\}$ defined by Algorithm 1 satisfies equation

$$
\mathcal{M} x_{k+1}=\mathcal{N} x_{k}+b
$$

Residual vectors of $x_{k+1}$ can be expressed as

$$
r_{k+1}=\left(\mathcal{I}-\mathcal{A} \mathcal{M}^{-1}\right) r_{k}=\mathcal{W} r_{k}
$$

and

$$
r_{k+1} \perp\left(\mathcal{I}-\frac{r_{k}^{T} \mathcal{W} r_{k}}{r_{k}^{T} \mathcal{W}^{T} \mathcal{W} r_{k}} \mathcal{W}\right) r_{k} \stackrel{\text { def }}{=} w_{k}
$$

The goal is to find such $\alpha$ that fulfills the following

$$
\underset{\alpha>0}{\arg \min }\left\|\mathcal{W} r_{k}\right\| /\left\|r_{k}\right\| .
$$

In other words we are looking for such $\alpha$ that the angle between vectors $r_{k}$ and $w_{k}$ is minimal. However, such formulated task leads on solution of nonlinear problem. Hence a uniform estimate for the equation $\left\|\mathcal{W} r_{k}\right\| /\left\|r_{k}\right\|$ is applied, i.e. $\alpha^{*}$ (optimal $\alpha$ ) is computed from the relation

$$
\begin{equation*}
\alpha^{*}=\underset{\alpha>0}{\arg \min }\left\|(\alpha \mathcal{I}-\mathcal{G})(\mathcal{G}+\alpha \mathcal{I})^{-1}\right\| \tag{12}
\end{equation*}
$$

Let us remark, that the matrix $(\alpha \mathcal{I}-\mathcal{S})(\alpha \mathcal{I}+\mathcal{S})^{-1}$ is orthonormal $((\Leftarrow \mathcal{S}$ is skewsymmetric) and can be therefore omitted in the norm in (12).

We continue to explore the expression (12): Under the assumption that $G$ is positive definite, $C$ is symmetric positive semidefinite and $B$ has maximum rank, it was shown in [7] that

$$
\begin{equation*}
\alpha^{*}=\underset{\alpha>0}{\arg } \min \max _{\substack{\lambda \in\left\{\lambda_{1}, \lambda_{2}, \ldots, \lambda_{n}\right\} \\ \lambda \in\left[\lambda_{\min }, \lambda_{\max }\right]}}\left|\frac{\alpha-\lambda}{\alpha+\lambda}\right|, \tag{13}
\end{equation*}
$$

where

$$
\lambda_{\min }:=\lambda_{1} \leq \lambda_{2} \leq \ldots \leq \lambda_{n}=: \lambda_{\max }
$$

are eigenvalues of the matrix $G$. The solution (13) yields

$$
\alpha *=\sqrt{\lambda_{\min } \lambda_{\max }} .
$$

Moreover for the condition number of the matrix we have the estimate

$$
\kappa\left(G+\alpha^{*} I\right) \leq 1+\sqrt{\kappa(G)} .
$$

where $G$ is the first diagonal block of matrix $\mathcal{G}$.
In the following will be shortly shown, how the relation (13) was developed and why the eigenvalues of matrix $C$ disappeared.

Let

$$
\mathcal{Z}=(\alpha \mathcal{I}-\mathcal{G})(\alpha \mathcal{I}+\mathcal{G})^{-1}(\alpha \mathcal{I}-\mathcal{S})(\alpha \mathcal{I}+\mathcal{S})^{-1}
$$

and

$$
\nu_{1} \leq \nu_{2} \leq \ldots \leq \nu_{m}
$$

be the eigenvalues of the matrix $C$. Because matrix $\mathcal{G}$ is symmetric, and matrix $(\alpha \mathcal{I}-\mathcal{S})(\alpha \mathcal{I}+\mathcal{S})^{-1}$ is orthonormal, there exists an orthonormal matrix $\mathcal{P}$ such, that

$$
\mathcal{P}^{T} \mathcal{Z P}=\underbrace{\mathcal{P}^{T}(\alpha \mathcal{I}-\mathcal{G})(\alpha \mathcal{I}+\mathcal{G})^{-1} \mathcal{P}}_{\left[\begin{array}{cc}
D_{1} & 0  \tag{14}\\
0 & D_{2}
\end{array}\right]=: \mathcal{D}} \underbrace{\mathcal{Q}}_{\mathcal{P}^{T}(\alpha \mathcal{I}-\mathcal{S})(\alpha \mathcal{I}+\mathcal{S})^{-1} \mathcal{P}},
$$

where

$$
\mathcal{Q}:=\left[\begin{array}{ll}
Q_{11} & Q_{12} \\
Q_{21} & Q_{22}
\end{array}\right]
$$

is orthonormal and

$$
D_{1}=\operatorname{diag}\left(\frac{\alpha-\lambda_{1}}{\alpha+\lambda_{1}}, \ldots \frac{\alpha-\lambda_{n}}{\alpha+\lambda_{n}}\right), \quad \varrho\left(D_{1}\right)<1
$$

and analogically with $D_{2}$, where we formally substitute $\nu_{i}$ for $\lambda_{i}, i=1, \ldots, m$. It holds $\varrho\left(D_{2}\right) \leq 1$. We try to show, that $\varrho(\mathcal{Q D})<1$.

Let $\left(\lambda,\left[x_{1}, x_{2}\right]^{T}\right)$ be an eigenpair of $\mathcal{Q D}$, that means

$$
\left[\begin{array}{ll}
Q_{11} D_{1} & Q_{12} D_{2} \\
Q_{21} D_{1} & Q_{22} D_{2}
\end{array}\right]\left[\begin{array}{l}
x_{1} \\
x_{2}
\end{array}\right]=\lambda\left[\begin{array}{l}
x_{1} \\
x_{2}
\end{array}\right] \stackrel{\text { if }}{\underline{x_{1}}=0}\left[\begin{array}{l}
Q_{12} D_{2} x_{2} \\
Q_{22} D_{2} x_{2}
\end{array}\right] \text { if } \stackrel{x_{1}=0}{=}\left[\begin{array}{c}
0 \\
\lambda x_{2}
\end{array}\right] .
$$

It remains to show that for all eigenvalues $\lambda$ the corresponding subvector $x_{1}$ is not a nullvector. From the last equation we immediately obtain

$$
\begin{equation*}
|\lambda|^{2}=\left\|D_{1} x_{1}\right\|^{2}+\left\|D_{2} x_{2}\right\|^{2} \leq \varrho\left(D_{1}\right)^{2}\left\|x_{1}\right\|^{2}+\varrho\left(D_{1}\right)^{2}\left\|x_{2}\right\|^{2}<\|x\|^{2}=1 . \tag{15}
\end{equation*}
$$

In the rest of our proof we will show that the matrix $Q_{12}$ has maximum rank. If we denote

$$
(\alpha \mathcal{I}-\mathcal{S})(\alpha \mathcal{I}+\mathcal{S})^{-1}=\underbrace{\left[\begin{array}{ll}
U_{11} & U_{12} \\
U_{21} & U_{22}
\end{array}\right]}_{\mathcal{U}}
$$

than

$$
\begin{equation*}
U_{12}=-2 \alpha(\alpha I+S) B^{T}\left[\alpha I_{m}+B\left(\alpha I_{m}+S\right)^{-1} B^{T}\right]^{-1} \tag{16}
\end{equation*}
$$

Using (4) we evaluate $\mathcal{Q}$ from (14):

$$
\mathcal{Q}=\underbrace{\left[\begin{array}{cc}
P_{11} & 0 \\
0 & P_{22}
\end{array}\right]^{T}}_{\mathcal{P}^{T}} \underbrace{\left[\begin{array}{ll}
U_{11} & U_{12} \\
U_{21} & U_{22}
\end{array}\right]}_{\mathcal{U}} \underbrace{\left[\begin{array}{cc}
P_{11} & 0 \\
0 & P_{22}
\end{array}\right]}_{\mathcal{P}}=\left[\begin{array}{cc}
P_{11}^{T} U_{11} P_{11} & P_{11}^{T} U_{12} P_{22} \\
P_{22}^{T} U_{21} P_{11} & P_{22}^{T} U_{22} P_{22}
\end{array}\right] .
$$

Hence

$$
Q_{12}=P_{11}^{T} U_{12} P_{22}=-2 \alpha P_{11}^{T}(\alpha I+S) B^{T} \underbrace{\left[\alpha I_{m}+B\left(\alpha I_{m}+S\right)^{-1} B^{T}\right]^{-1}}_{E} P_{22}
$$

where we have substituted from (16). Matrices $P_{11}$ and $P_{22}$ are orthonormal, the matrix $E$ is nonsingular and $B$ has maximal column rank, hence the matrix $Q_{12}$ has maximum rank.

## 3. Polynomial filters

Assume that the eigenvalues of the matrix $G$ are ordered according to

$$
0<\lambda_{1} \leq \lambda_{2} \leq \ldots \leq \lambda_{n}
$$

Let $k \ll n$ be a positive integer. We want to find the eigenvalues included in the set $\sigma_{w}=\left\{\lambda_{1}, \ldots, \lambda_{k}\right\}$ and $\sigma_{u}=\left\{\lambda_{n-k+1}, \ldots, \lambda_{n}\right\}$ respectively.
Lanczos process starting with the vector $v_{1}$ gives the factorization

$$
\begin{equation*}
G V_{m}=V_{m} T_{m}+f_{m} e_{m}^{T} \tag{17}
\end{equation*}
$$

where $V_{m}=\left\{v_{1}, v_{2}, \ldots, v_{m}\right\} \in \mathbb{R}^{n \times m}, f_{m} \in \mathbb{R}^{n}, T_{m} \in \mathbb{R}^{m \times m}, V_{m}^{T} V_{m}=I_{m}, V_{m}^{T} f_{m}=0$.

We will consider two ways for estimating the eigenvalues in $\sigma_{w}$ or $\sigma_{u}$. Either for larger $m$ to calculate the smallest and largest eigenvalues of the tridiagonal matrix $T_{m}$ or for small $m$ repeat the procedure based on the single-shift QR iterations to obtain more detail information about spectrum of the matrix $G$. The second procedure is analogous to the updating the Arnoldi factorization via QR-iteration (see References [24], [11]) and we describe it shortly. Let $\mu_{1}, \mu_{2}, \ldots, \mu_{m-k}$ be positive real numbers. From (17) we successively obtain

$$
\begin{align*}
\left(G-\mu_{1} I\right) V_{m}-V_{m}\left(T_{m}-\mu_{1} I\right) & =f_{m} e_{m}^{T},  \tag{18}\\
\left(G-\mu_{1} I\right) V_{m}-V_{m} Q^{(1)} R^{(1)} & =f_{m} e_{m}^{T},  \tag{19}\\
\left(G-\mu_{1} I\right)\left(V_{m} Q^{(1)}\right)-\left(V_{m} Q^{(1)}\right)\left(R^{(1)} Q^{(1)}\right) & =f_{m} e_{m}^{T} Q^{(1)},  \tag{20}\\
G \underbrace{\left(V_{m} Q^{(1)}\right)}_{V_{m}^{(1)}}-\underbrace{\left(V_{m} Q^{(1)}\right)}_{V_{m}^{(1)}} \underbrace{\left(R^{(1)} Q^{(1)}+\mu_{1} I\right)}_{T_{m}^{(1)}} & =f_{m} e_{m}^{T} Q^{(1)} . \tag{21}
\end{align*}
$$

Analogical application of the next shifts $\mu_{2}, \mu_{3}, \ldots, \mu_{m-k}$ gives

$$
\begin{equation*}
G V_{m}^{(m-k)}=V_{m}^{(m-k)} T_{m}^{(m-k)}+f_{m} e_{m}^{T} Q, \tag{22}
\end{equation*}
$$

where

$$
\begin{aligned}
& Q=Q^{(1)} Q^{(2)} \ldots Q^{(m-k)}, \quad T_{m}^{(m-k)}=Q^{T} T_{m} Q \quad \text { and } \\
& V_{m}^{(m-k)} ;=\left(v_{1}^{(m-k)}, v_{2}^{(m-k)}, \ldots, v_{m}^{(m-k)}\right)=V_{m} Q
\end{aligned}
$$

The matrix $T_{m}^{(m-k)}$ has the form

$$
T_{m}^{(m-k)}=\left(\begin{array}{cc}
T_{k}^{(m-k)} & t_{k} e_{k} e_{1}^{T} \\
t_{k} e_{1} e_{k}^{T} & \widehat{T}_{m-k}
\end{array}\right),
$$

and if we split $V_{m}^{(m-k)}=\left(V_{k}^{(m-k)}, \tilde{V}_{m-k}^{(m-k)}\right)$ then

$$
\begin{equation*}
G V_{k}^{(m-k)}=V_{k}^{(m-k)} T_{k}^{(m-k)}+\underbrace{f_{k}^{(m-k)}}_{t_{k+1, k} v_{k+1}^{(m-k)}+f_{m} e_{m}^{T} Q e_{k}^{T}} e_{k}^{T} \tag{23}
\end{equation*}
$$

where $V_{k}^{T} f_{k}^{(m-k)}=0$ and the matrices $T_{m}^{(m-k)}$ and $T_{k}^{(m-k)}$ are tridiagonal positive definite because $G$ is positive definite. The equalities (18)-(21) yield the formula

$$
\begin{equation*}
v_{1}^{m-k}=\nu \Pi_{j=1}^{m-k}\left(A-\mu_{j} I\right) v_{1} \stackrel{\text { def }}{=} \psi(A) v_{1} \tag{24}
\end{equation*}
$$

The equation (23) gives Lanczos process with starting vector $V_{k}^{(m-k)} e_{1}$. Let

$$
\begin{equation*}
0<\theta_{1} \leq \theta_{2} \leq \ldots \leq \theta_{m} \tag{25}
\end{equation*}
$$

be the eigenvalues of the matrix $T_{m}$. Now it is natural to put the following question. How to choose the shifts $\mu_{j}$ or alternatively a polynomial $\psi$ such that $\left\|f_{k}^{m-k}\right\|$ is
"small", because if it would be in ideal case $\left\|f_{k}^{m-k}\right\|=0$ then the columns of the matrix $V_{k}^{m-k}$ form an invariant subspace of the matrix $G$. Let tol is a small positive number (for example tol $=10^{-2}$ ). This will be discussed in the section concerning with numerical results. Let $v_{1}$ be the vector defined in (17). The polynomial $\psi$ having the property that the Lanczos process starting with the initial vector $\psi(G) v_{1}$ gives $\left\|f_{k}^{m-k}\right\| \leq t o l$ will be called polynomial filter. If the polynomial $\psi$ is constructed iteratively then all polynomial approximations will be called filters.

However we want to get eigenspace corresponding to the smallest or largest eigenvalues of $G$. Let us think the smallest eigenvalues. The algorithm for largest eigenvalues can be implemented equally. By analogy to the [24] paper let us consider the following iterative process.

```
Algorithm 2: Construction of an invariant subspace
Input: \(G, m, k, t o l\), max - maximal number of iterations
Output: matrices \(V_{k}\) and a tridiagonal matrix \(T_{k}\) fulfilling the relations
    \(\left\|G V_{k}-V_{k} T_{k}\right\| \leq t o l, V_{k}^{T} V_{k}=I\)
Step 1 Carry out the Lanczos process (17), cycle \(:=1\)
Step 2 Calculate the eigenvalues \(\theta_{j}, j=1,2, \ldots, m\) of \(T_{m}\)
Step 3 Substitute \(\mu_{j}:=\theta_{m-j}\) for \(j=1,2, \ldots, k\) and apply these shifts to
    obtain (22) and (23).
Step 4 If \(\left\|f_{k}^{m-k}\right\|<\) tol or cycle \(>\) max then STOP else
    perform the next \(m-k\) steps of the the Lanczos process starting
    from (23). The resulting tridiagonal matrix is denoted again \(T_{m}\).
    Put cycle \(=\) cycle +1 , go to Step 2 .
end.
```

The convergence is studied in [24]. Let us assume that the process converges, i.e. $\left\|f_{k}^{(m-k)}\right\|<t o l$ after performing a finite number of steps.
We have

$$
A V_{k} \doteq V_{k} T_{k}
$$

and $\operatorname{span}\left\{V_{k}\right\}$ approximate eigenspace corresponding to the $k$ eigenvalues of $\lambda_{1}$, $\lambda_{2}, \ldots \lambda_{k} G$. (See Reference [24].) If $\left\{\theta_{j}^{(k)}, y_{j}^{(k)}\right\}_{j=1}^{k}$ are eigenpairs of $T_{k}$, and

$$
x_{j}=V_{k} y_{j}^{(k)}, \quad j=1,2, \ldots, k \text { the Ritz vectors }
$$

then the easy manipulation shows that the pairs $\left(\theta_{j}^{(k)}, x_{j}\right)$ approximate the eigenpairs of $G$ and, moreover, the eigenpairs $\left(\lambda_{j}, u_{j}\right)$, where the matrix $U=\left(u_{1}, u_{2}, \ldots, u_{n}\right)$ transforms $G$ to the Jordan canonical form $\Lambda=\operatorname{diag}\left(\lambda_{1}, \lambda_{2}, \ldots, \lambda_{n}\right)$. It is an interesting property proved by [24] that the $k$ smallest Ritz eigenvalues (i.e. smallest eigenvalues of the matrix $T_{m}$ ) approximate successively $\left(\lambda_{1}, \lambda_{2}, \ldots, \lambda_{k}\right)$. (See Reference [24])

## 4. The Chebyshev filtering technique

Let

$$
\begin{equation*}
G=U \Lambda U^{T}=U_{1} \Lambda_{1} U_{1}^{T}+U_{2} \Lambda_{2} U_{2}^{T} \tag{26}
\end{equation*}
$$

where $\Lambda_{1}=\operatorname{diag}\left(\lambda_{1}, \lambda_{2}, \ldots, \lambda_{k}\right)$ and $\Lambda_{2}=\operatorname{diag}\left(\lambda_{k+1}, \lambda_{2}, \ldots, \lambda_{n}\right.$ be the eigendecomposition of the symmetric and positive definite matrix $G$. Let $v_{1}$ be a starting vector for Lanczos process and $\psi$ an arbitrary polynomial. From (26)it follows that

$$
\begin{equation*}
\psi(G) v_{1}=U_{1} \psi\left(\Lambda_{1}\right) U_{1}^{T} v_{1}+U_{2} \psi\left(\Lambda_{2}\right) U_{2}^{T} v_{1} \tag{27}
\end{equation*}
$$

We want to find $\psi$ which sets off the first term on the right hand-side (27) and damps the second one. The min-max property of Chebyshev polynomials leads us to the idea to seek a polynomial $\psi$ between them. Generally, the Chebyshev polynomials $T_{m}$ of degree $m$ are defined for all complex numbers by the formula $T_{m}(z)=\cosh \left(m \times \cosh ^{-1}(z)\right)$ as an analytical function (see Reference [25]). But all eigenvalues of $G$ are real and hence real version of these polynomials will be sufficient for our further investigations. For the interval $[-1,1]$ the polynomials $T_{m}$ reduce to very easy form $T_{m}(z)=\cos \left(m \times \cos ^{-1}(z)\right)$. The practical algorithms utilize often the recurrence formulas

$$
T_{0}(z)=1, \quad T_{1}(z)=z \quad \text { and } \quad T_{s+1}(z)=2 z T_{s}(z)-T_{s-1}(z) \quad \text { for all integer } s \geq 1
$$

We want theoretically "annihilate" the norm of the matrix $\psi\left(\Lambda_{2}\right)$. for the approximation of the smallest interval containing diagonal elements of $\Lambda_{2}$ we apply again Lanczos process and Ritz values. Let $m$ and $k$ has the same meaning as in the formulas (17) and (22). Let $\gamma$ and $\delta$ be a positive numbers. The transformation

$$
\begin{equation*}
z \mapsto \vartheta(\lambda):=\frac{\delta+\gamma-2 \lambda}{\delta-\gamma} \tag{28}
\end{equation*}
$$

maps the interval $[\gamma, \delta]$ onto $[-1,1]$ and $\vartheta(0)=\frac{\delta+\gamma}{\delta-\gamma}=: q$. Hence assuming the matrix $G$ possesses the unwanted (the largest or smallest) eigenvalues in the interval $[\gamma, \delta]$, the polynomial $\psi$,

$$
\psi(\lambda)=\frac{T_{m-k}(\vartheta(\lambda))}{T_{m-k}(q)}
$$

normed that $\left\|\psi(G) v_{1}\right\|=1$ appears to be a good polynomial filter. Let us remark that an analogous technique was used in the papers [Arioli, Ruiz atd.] where the Chebyshev damping with the block Lanczos method is effectively used for solving large linear systems with a symmetric and positive definite matrix. Now the question appears how numerically perform this damping of eigenvalues lying in $[\gamma, \delta]$ and how to apply the product $\psi(A) v_{1}$. The process is very similar to the Algorithm 1 and we underline here only the differences. Denote $t=m-k$ and carry out the Lanczos process (17), i.e. the Step 1 in Algorithm 1 and we calculate all eigenvalues of $T_{m}$, see Step 2.

We put $\gamma \in\left(0, \theta_{1}\right]$ and $\delta \in\left[\theta_{k} \theta_{k+1)}\right.$ ) if small eigenvalues are damped, i.e. the invariant subspace for $t$ largest eigenvalues of $G$ is calculated or we substitute $\gamma \in$ $\left(\theta_{k-1} \theta_{k}\right]$ and $\delta \in\left[\theta_{m}, \chi\right)$ for some $\chi>\theta_{m}$ in the opposite case. Let us still remark that the linear mapping

$$
\begin{equation*}
\lambda \mapsto \eta(\theta):=\delta+(1-\theta)(\gamma-\delta) / 2 \tag{29}
\end{equation*}
$$

transforms the interval $[-1,1]$ onto $[\gamma, \delta]$ and to the roots of $T_{m}$ in $[-1,1]$ correspond the numbers $\mu_{1}, \mu_{2}, \ldots \mu_{r}$, where $r=k$ or $r=t$ dependent if the invariant subspace for $t$ largest eigenvalues of $G$ is calculated or conversely. Hence for the shifts we take the numbers $\mu_{j}$

$$
\mu_{j}=\eta\left(\cos \left(\frac{2 j-1}{2 t} \pi\right)\right) \quad \text { for } \quad j=1,2, \ldots r .
$$

The Step 4 stays only instead of $m-k$ must be substituted $k$ if largest eigenvalues are calculated.

## 5. Solution of systems with skew-symmetric matrix

In this section we deal with systems of the form

$$
\begin{aligned}
(\alpha \mathcal{I}+\mathcal{S}) x^{k+1} & =(\alpha \mathcal{I}-\mathcal{G}) x^{k+\frac{1}{2}}+b, \quad \text { and } \\
(\underbrace{\left(\left(\alpha I_{n}+\frac{1}{\alpha} B^{T} B\right)\right.}_{Q} & +S) u^{k+1}=f_{2}^{k}-\frac{1}{\alpha} B^{T} g_{2}^{k}
\end{aligned}
$$

The matrices of both systems have the form: "a SPD matrix + an skew-symmetric one" and therefore without any loss of generality we will consider only the second equation, we leave out the upper index $k$ and denote the-right hand side by $f_{r}$, i.e. we mean the system

$$
\begin{equation*}
(Q+S) u=f_{r} \tag{30}
\end{equation*}
$$

Let us mention that $Q$ is symmetric and positive definite and $S$ is skew-symmetric matrix. The acceleration procedure for the CGW iterative method, formulated by Concus, Golub and Widlund, is recommended (see References [22], [7]). It accelerates the iterative process

$$
\underbrace{c=Q^{-1} b .}_{u^{i+1}=T u^{i}+c \quad \underbrace{\left.u^{i+1}=Q^{-1}(Q+S)\right)}_{-Q^{-1} S=T} u^{i}+Q^{-1} b}
$$

Let us shortly present following iteration formulas:

$$
\begin{aligned}
u^{i+1} & =\omega_{i+1}\left(T u^{i}+c\right)+\left(1-\omega_{i+1}\right) u^{i-1} \\
\omega_{i+1} & =\left(1-\frac{\left(Q \delta^{i}\right)^{T} \delta^{i}}{\left(Q \delta^{i-1}\right)^{T} \delta^{i-1}} \frac{1}{\omega_{1}}\right)^{-1} \omega_{1}=1 .
\end{aligned}
$$

where

$$
\delta^{i}=T u^{i}+c-u^{i} .
$$

Let us remark, that the last formulas may be fast obtained by rewriting of the three-term recurrence formula for the preconditioned CG applied in symmetric case.

## 6. Numerical results

In this section, we compare some Krylov subspace methods with the alternating iteration method described in this paper. We show, that this method is faster for larger systems of the considered structure.

As it was shown, one of our main goals is to find an optimal parameter $\alpha^{*}$ to gain fast convergence of Algorithm 1. Therefore we have to calculate the smallest and largest eigenvalue of matrix $G$. We test the Lanczos method and the method described in Algorithm 2. In the second case we obtain more than one of the smallest or largest eigenvalues and this fact gives an information, which is used by the construction of a good preconditioner for the system (G).

For the solution of the symmetric part in Algorithm 1 we use preconditioned CG method. The preconditioning matrix $M$ was constructed according to the algorithm presented in [23], where $M$ was constructed analogously to the paper [11]. This preconditioner is very efficient in the case of matrices having very small eigenvalues. The first example will be academic.
Example. Let $A \in \mathbb{R}^{100 \times 100}$ be symmetric positive definite matrix with eigenvalues $0.009887,0.01803,0.03207$. The other ones lie in interval [3.254, 100.7]. Hence our matrix $A$ has three very small eigenvalues. We use the Algorithm 2 with parameters $s=8, t=3$ and after 29 iterations we get the three smallest eigenvalues of $A$ with precision $10^{-14}$ and the corresponding eigenvectors. If we solve linear system $A z=e$ using preconditioned CG method (for the construction of preconditioner were used polynomial filters). The residual norm is less then $10^{-4}$ after 140 iterations. Classical CG-method or CG-method preconditioned with incomplete Cholesky stagnates.

### 6.1. Test cases

Our concern is to solve large linear systems occurring in fluid dynamics, mostly constructed by finite element or finite volume method.

Let us consider a flow in the unit square domain $\Omega$ described by the following equations

$$
\begin{aligned}
& -\mu \Delta u_{1}+\left\{v_{1} \cdot\left(u_{1}\right)_{x}+v_{2} \cdot\left(u_{1}\right)_{y}\right\}+p_{x}=0 \quad \text { in } \Omega \\
& -\mu \triangle u_{2}+\left\{v_{1} \cdot\left(u_{2}\right)_{x}+v_{2}\left(u_{2}\right)_{y}\right\}+p_{y}=0 \text { in } \Omega \\
& \left(u_{1}\right)_{x}+\left(u_{2}\right)_{y}=0 \text { in } \Omega \\
& u_{1}=1 \text { on } \Gamma_{1} \\
& u_{1}=0 \text { on } \Gamma \backslash \Gamma_{1} \\
& u_{2}=0 \text { on } \Gamma
\end{aligned}
$$

where $\Gamma$ is the boundary of the unit square and $\Gamma_{1}$ his part with $y=1$. The functions $u_{1}$ and $u_{2}$ are the velocity components in $x$ and $y$ directions and $p$ is the

| Name | Size |  | Parameters |
| :--- | ---: | ---: | :--- |
|  | Nr. of unknowns | Nr. of non-zero el. |  |
| Stokes 1 | 759 | 3764 | $h=0.1, \nu=10^{-1}$ |
| Stokes 2 | 3119 | 15924 | $h=0.05, \nu=10^{-1}$ |
| Stokes 3 | 7079 | 36484 | $h=0.033, \nu=10^{-1}$ |
| Stokes 4 | 1263 | 65444 | $h=0.025, \nu=10^{-1}$ |
| Stokes 5 | 19799 | 102804 | $h=0.02, \nu=10^{-1}$ |
| Stokes 6 | 28559 | 148564 | $h=0.0167, \nu=10^{-1}$ |
| Stokes 7 | 38919 | 202724 | $h=0.0143, \nu=10^{-1}$ |
| Stokes 8 | 50879 | 265284 | $h=0.0125, \nu=10^{-1}$ |
| Stokes 9 | 64439 | 336244 | $h=0.0111, \nu=10^{-1}$ |
| Stokes 10 | 79599 | 415604 | $h=0.01, \nu=10^{-1}$ |

Tab. 1: Set of test matrices.
pressure. The resulting linear system for the discrete solution, putting $u=\left[u_{1}, u_{2}\right]^{T}$, has the form (1) with $C=0$ and $A$ symmetric and positive definite.

### 6.2. Estimation of optimal parameter

We want to find an optimal parameter $\alpha^{*}=\sqrt{\lambda_{\min } \lambda_{\max }}$ for the matrices defined in Table 1. We compare the Lanczos method with the Algorithm 2. In Table 2 the following notation is used:

- $\varepsilon_{\min }=\left|\lambda_{\min }^{(\text {comp })}(G)-\lambda_{\min }(G)\right|$, where $\lambda_{\min }^{(\text {comp })}$ is computed and $\lambda_{\text {min }}$ is the exact smallest eigenvalue of $G$.
- $\varepsilon_{\max }=\left|\lambda_{\max }^{(\text {comp })}(G)-\lambda_{\max }(G)\right|$, where $\lambda_{\max }^{(\text {comp })}$ is computed and $\lambda_{\max }$ is the exact largest eigenvalue of $G$.
- fl. $=$ number of floating-point operations.
- $s=$ number of Lanczos steps.
- $s_{\text {opt }}=$ number of steps necessary to to obtain results comparable with Algorithm 2.
- $m, k=$ constants defined in Algorithm 2.

We show the results obtained using Laczos method for $s=10$ and $s_{\text {opt }}$ in the first two columns. We choose constant $s$, such that the computation costs are similar to Algorithm 2. The last two columns represents the Algorithm 2. In column 3 and 4 different variants of shifts are used according to the text in Section 3.

Table 2 illustrates, that for smaller problems the Lanczos method gives similar results as Algorithm 2. But when the size of the problem grows Algorithm 2 appears to be more efficient.

| Matrix | Lanczos |  | Chebyshev | Sorensen |
| :---: | :---: | :---: | :---: | :---: |
|  | $s=10$ | $s=o p t$ | $m=5, k=3$ | $m=5, k=3$ |
| Stokes 1 | $\begin{aligned} & \varepsilon_{\min }=10^{-3} \\ & \varepsilon_{\max }=10^{-2} \\ & \text { fl. }=186742 \end{aligned}$ | $\begin{aligned} & s_{\text {opt }}=60 \\ & \text { fl. }=1164243 \end{aligned}$ | $\begin{aligned} & \varepsilon_{\min }=10^{-5} \\ & \varepsilon_{\max }=10^{-4} \\ & \text { fl. }=152433 \end{aligned}$ | $\begin{aligned} & \varepsilon_{\min }=10^{-5} \\ & \varepsilon_{\max }=10^{-5} \\ & \mathrm{fl}=156476 \end{aligned}$ |
| Stokes 2 | $\begin{aligned} & \varepsilon_{\min }=10^{-1} \\ & \varepsilon_{\max }=10^{-1} \\ & \text { f. }=0.8 \times 10^{6} \end{aligned}$ | $\begin{aligned} & s_{\text {opt }}=65 \\ & \text { fl. }=1.9 \times 10^{7} \end{aligned}$ | $\begin{aligned} & \varepsilon_{\min }=10^{-5} \\ & \varepsilon_{\max }=10^{-5} \\ & \text { fl. }=8.86 \times 10^{6} \end{aligned}$ | $\begin{aligned} & \varepsilon_{\min }=10^{-5} \\ & \varepsilon_{\max }=10^{-5} \\ & \text { fl. }=8.87 \times 10^{6} \end{aligned}$ |
| Stokes 3 | $\begin{aligned} & \varepsilon_{\min }=10^{-1} \\ & \varepsilon_{\max }=10^{-1} \\ & \text { f. }=1.8 \times 10^{6} \end{aligned}$ | $\begin{aligned} & s_{\text {opt }}=68 \\ & \mathrm{fl}=2.4 \times 10^{7} \end{aligned}$ | $\begin{aligned} & \varepsilon_{\min }=10^{-5} \\ & \varepsilon_{\max }=10^{-4} \\ & \text { fl. }=8.91 \times 10^{6} \end{aligned}$ | $\begin{aligned} & \varepsilon_{\min }=10^{-6} \\ & \varepsilon_{\max }=10^{-6} \\ & \text { fl. }=8.93 \times 10^{6} \end{aligned}$ |
| Stokes 4 | $\begin{aligned} & \varepsilon_{\min }=10^{-1} \\ & \varepsilon_{\max }=10^{-1} \\ & \text { f. }=3.2 \times 10^{6} \end{aligned}$ | $\begin{aligned} & s_{\text {opt }}=74 \\ & \text { fl. }=7.4 \times 10^{7} \end{aligned}$ | $\begin{aligned} & \varepsilon_{\min }=10^{-5} \\ & \varepsilon_{\max }=10^{-5} \\ & \text { fl. }=9.67 \times 10^{6} \end{aligned}$ | $\begin{aligned} & \varepsilon_{\min }=10^{-7} \\ & \varepsilon_{\max }=10^{-5} \\ & \text { fl. }=9.98 \times 10^{6} \end{aligned}$ |
| Stokes 5 | $\begin{aligned} & \varepsilon_{\min }=10^{-1} \\ & \varepsilon_{\max }=10^{-1} \\ & \mathrm{fl} .=5.0 \times 10^{6} \end{aligned}$ | $\begin{aligned} & s_{\text {opt }}=74 \\ & \text { fl. }=8.4 \times 10^{7} \end{aligned}$ | $\begin{aligned} & \varepsilon_{\min }=10^{-6} \\ & \varepsilon_{\max }=10^{-6} \\ & \mathrm{fl} .=1.02 \times 10^{7} \end{aligned}$ | $\begin{aligned} & \varepsilon_{\min }=10^{-5} \\ & \varepsilon_{\max }=10^{-5} \\ & \text { fl. }=1.02 \times 10^{7} \end{aligned}$ |
| Stokes 6 | $\begin{aligned} & \varepsilon_{\min }=10^{-1} \\ & \varepsilon_{\max }=10^{-1} \\ & \mathrm{fl} .=7.2 \times 10^{6} \end{aligned}$ | $\begin{aligned} & s_{\text {opt }}=76 \\ & \text { fl. }=1.2 \times 10^{8} \end{aligned}$ | $\begin{aligned} & \varepsilon_{\min }=10^{-5} \\ & \varepsilon_{\max }=10^{-5} \\ & \text { fl. }=1.11 \times 10^{7} \end{aligned}$ | $\begin{aligned} & \varepsilon_{\min }=10^{-5} \\ & \varepsilon_{\max }=10^{-5} \\ & \text { fl. }=1.12 \times 10^{7} \end{aligned}$ |
| Stokes 7 | $\begin{aligned} & \varepsilon_{\min }=10^{-1} \\ & \varepsilon_{\max }=10^{-1} \\ & \text { fl. }=9.8 \times 10^{6} \end{aligned}$ | $\begin{aligned} & s_{\text {opt }}=78 \\ & \text { fl. }=1.6 \times 10^{8} \end{aligned}$ | $\begin{aligned} & \varepsilon_{\min }=10^{-5} \\ & \varepsilon_{\max }=10^{-5} \\ & \text { fl. }=1.32 \times 10^{7} \end{aligned}$ | $\begin{aligned} & \varepsilon_{\min }=10^{-5} \\ & \varepsilon_{\max }=10^{-5} \\ & \text { fl. }=1.35 \times 10^{7} \end{aligned}$ |
| Stokes 8 | $\begin{aligned} & \varepsilon_{\min }=10^{-1} \\ & \varepsilon_{\max }=10^{-1} \\ & \text { f. }=1.3 \times 10^{7} \end{aligned}$ | $\begin{aligned} & s_{\text {opt }}=76 \\ & \text { fl. }=2.2 \times 10^{8} \end{aligned}$ | $\begin{aligned} & \varepsilon_{\min }=10^{-5} \\ & \varepsilon_{\max }=10^{-5} \\ & \text { fl. }=1.43 \times 10^{7} \end{aligned}$ | $\begin{aligned} & \varepsilon_{\min }=10^{-5} \\ & \varepsilon_{\max }=10^{-5} \\ & \text { fl. }=1.44 \times 10^{7} \end{aligned}$ |
| Stokes 9 | $\begin{aligned} & \varepsilon_{\min }=10^{-1} \\ & \varepsilon_{\max }=10^{-1} \\ & \text { f. }=1.6 \times 10^{7} \end{aligned}$ | $\begin{aligned} & s_{\text {opt }}=84 \\ & \text { fl. }=2.6 \times 10^{8} \end{aligned}$ | $\begin{aligned} & \varepsilon_{\min }=10^{-5} \\ & \varepsilon_{\max }=10^{-5} \\ & \text { fl. }=1.56 \times 10^{7} \end{aligned}$ | $\begin{aligned} & \varepsilon_{\min }=10^{-5} \\ & \varepsilon_{\max }=10^{-5} \\ & \text { fl. }=1.57 \times 10^{7} \end{aligned}$ |
| Stokes 10 | $\begin{aligned} & \varepsilon_{\min }=10^{-1} \\ & \varepsilon_{\max }=10^{-1} \\ & \text { f. }=1.9 \times 10^{7} \end{aligned}$ | $\begin{aligned} & s_{\text {opt }}=89 \\ & \text { fl. }=2.9 \times 10^{8} \end{aligned}$ | $\begin{aligned} & \varepsilon_{\min }=10^{-5} \\ & \varepsilon_{\max }=10^{-5} \\ & \mathrm{fl}=1.65 \times 10^{7} \end{aligned}$ | $\begin{aligned} & \varepsilon_{\min }=10^{-5} \\ & \varepsilon_{\max }=10^{-5} \\ & \mathrm{fl}=1.66 \times 10^{7} \end{aligned}$ |

Tab. 2: Estimation of optimal parameter $\alpha^{*}$.

### 6.3. Methods comparison

In this subsection we compare GMRES, GMRES(20), MINRES and ADI methods for the solution of our test cases. In the following table we compare the floatingpoint operations necessary to get the solution with residual norm less than $10^{-6}$. The parameter $\alpha$ used in algorithm 1 was computed by Algorithm 2. The symbol $*$ means, that no results has been obtained in real time.

| Matrix | GMRES | GMRES(20) | MINRES | ADI |
| :--- | :--- | :--- | :--- | :---: |
| Stokes 1 | $1.1 \times 10^{6}$ | $0.5 \times 10^{6}$ | $0.7 \times 10^{6}$ | $2.1 \times 10^{6}$ |
| Stokes 2 | $4.2 \times 10^{6}$ | $1.2 \times 10^{6}$ | $1.1 \times 10^{6}$ | $3.2 \times 10^{6}$ |
| Stokes 3 | $8.4 \times 10^{6}$ | $2.1 \times 10^{6}$ | $1.9 \times 10^{6}$ | $4.4 \times 10^{6}$ |
| Stokes 4 | $*$ | $3.3 \times 10^{6}$ | $3.5 \times 10^{6}$ | $5.1 \times 10^{6}$ |
| Stokes 5 | $*$ | $6.1 \times 10^{6}$ | $5.6 \times 10^{6}$ | $7.3 \times 10^{6}$ |
| Stokes 6 | $*$ | $7.4 \times 10^{6}$ | $6.3 \times 10^{6}$ | $7.6 \times 10^{6}$ |
| Stokes 7 | $*$ | $8.2 \times 10^{6}$ | $*$ | $8.3 \times 10^{6}$ |
| Stokes 8 | $*$ | $9.3 \times 10^{6}$ | $*$ | $8.9 \times 10^{6}$ |
| Stokes 9 | $*$ | $1.6 \times 10^{7}$ | $*$ | $1.2 \times 10^{6}$ |
| Stokes 10 | $*$ | $2.8 \times 10^{7}$ | $*$ | $2.1 \times 10^{6}$ |

Tab. 3: Methods comparison.
We can see from Table 3, that the ADI method is more efficient if the size of the solved problem grows. Beginning with Stokes 7 the ADI method is better than GMRES(20).

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