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# OPTIMIZING TWO-LEVEL PRECONDITIONINGS FOR THE CONJUGATE GRADIENT METHOD 

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# Optimizing Two-Level Preconditionings for the Conjugate Gradient Method 

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#### Abstract

The construction of efficient iterative linear equation solvers for ill-conditioned general symmetric positive definite systems is discussed. Certain known twolevel conjugate gradient preconditioning techniques are presented in a uniform way and are further generalized and optimized with respect to the spectral or the K-condition numbers. The resulting constructions have shown to be useful for the solution of large-scale ill-conditioned symmetric positive definite linear systems.


## 1 Introduction

In the present paper, we address the construction of preconditionings for the Conjugate Gradient algorithm, see, e.g.[1], and the rate of convergence of the method. This method is used for solving linear algebraic systems

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

with a large, normally sparse, unstructured Symmetric Positive Definite (SPD) matrix $A$ of order $n$, such as arising in computational mechanics, from symmetrization of unsymmetric problems, etc.

Below we consider a preconditioning which is closely related to both the Generalized Augmented Matrix (GAM) preconditioning [18] and the approximate Schur complement one [4], [1], [3]. We restrict our considerations to two-level schemes based on a $2 \times 2$ splitting of the coefficient matrix and present a uniform framework for the analysis of such preconditioners. One of the main results is that the K-condition number of the preconditioned matrix is minimized under a very simple choice of approximation of the Schur complement. The upper bounds obtained for the Kcondition and spectral condition numbers of the preconditioned matrix show that one

[^0]can expect good overall preconditioning quality whenever the preconditioning of the leading block of the matrix has a sufficiently high quality. The latter can be attained by a proper choice of the $2 \times 2$ splitting of the matrix, as well as by application of improved preconditioning methods such as Second Order Cholesky type incomplete factorizations [16]. As shown in [8], for finite element applications of second order problems using a certain element based preconditioner, it is also possible to obtain accurate bounds of the condition number which hold uniformly in both problem and discretization parameters.

In the present paper we consider a combined preconditioning strategy for the Conjugate Gradient algorithm intended to achieve fast convergence while providing low iteration costs. The algorithm can be considered as a proper combination of preconditioning strategies described in $[4,7,18,12,13,15,16]$. In the first stage, the original matrix $A$ is split into a 2 by 2 block form with the leading block as large as possible while still well-conditioned. This is typically accompanied by a certain congruence transformation which is intended to enable a further improvement of the conditioning of the whole matrix or its leading block. In the second stage, an (approximate) block Jacobi preconditioning is used to construct the preconditioner in its final form.

The remainder of the paper is organized as follows. In Section 2 we recall two upper bounds on the number of iterations for the conjugate gradient method, and indicate their usefulness in the construction of preconditioners. In Section 3 a uniform presentation of various two-stage preconditionings is presented with condition number optimality results. In the same framework, a treatment of Schur complement preconditioners is given in Section 4. In Section 5 we describe the Conjugate Gradient Normal Equations algorithm for a guaranteed precision iterative solution of highly unsymmetric systems and discuss the potential use of the above described preconditionings with this method.

## 2 Two iteration bounds for the Preconditioned CG method

In order to highlight the target functions that should be optimized by the preconditioning, let us recall some known convergence results for the PCG method.

Consider the PCG method for the solution of SPD systems with symmetric positive definite preconditioner $H$ that approximates $A^{-1}$ in some sense. The standard estimation for the PCG iteration number needed for an $\varepsilon$ times reduction of the error norm $\left(r_{i}^{T} A^{-1} r_{i}\right)^{1 / 2}$ is (see. e.g., [1])

$$
i_{\kappa}(\varepsilon) \leq\left\lceil\frac{1}{2} \sqrt{\kappa(H A)} \log \frac{2}{\varepsilon}\right\rceil
$$

where, for any symmetrizable matrix $M$ with positive eigenvalues,

$$
\kappa(M)=\lambda_{\max }(M) / \lambda_{\min }(M)
$$

This bound follows from a well-known estimate, cf.[1], establishing the linear rate of convergence for the PCG iterations. In some (model) cases, the latter estimate is useful for obtaining a priori bounds expressed via the parameters of the problem solved. For an important example, see [2]. However, the requirement of "optimal" conditioning does not in general yield a concrete construction of the preconditioning.

Therefore, an alternative approach was developed based on the use of an iteration number estimate via the K-condition number, cf. $[15,1]$. Based on the corresponding superlinear convergence rate result, a simplified iteration number estimate of the following form holds (provided that the $A$-norm of the residual is replaced by the $H$-norm):

$$
i_{K}(\varepsilon)=\left\lceil\log _{2} K(H A)+\log _{2} \frac{1}{\varepsilon}\right\rceil
$$

where, by definition,

$$
K(M)=\left(\frac{1}{n} \operatorname{trace}(M)\right)^{n} / \operatorname{det}(M)
$$

This bound can be useful in predicting the superlinear rate of convergence for a number of iterations exceeding, but close to, $\log _{2} K(H A)$. It follows that

$$
i_{K}(\varepsilon) \leq\left\lceil n \log _{2}\left(\frac{a}{g}\right)+\log _{2}\left(\frac{1}{\varepsilon}\right)\right\rceil
$$

where $a$ is the arithmetic average and $g$ is the geometric average of the eigenvalues of $H A$. Typically, in practice for instance when considering a class of problems of increasing sizes, such as for difference methods for partial differential equations, it holds that $a / g \geq 1+c$ for some postive $c$, independent on $n$. Hence in such cases $i_{K}(\varepsilon) \leq c n+\log _{2} \frac{1}{\varepsilon}$.
Therefore this condition number somestimes gives rather pessimistic a priori upper bounds of the number of iterations. However, it may readily be (nearly) minimized in the context of various preconditioning procedures, as shown already in $[12,15,1]$. Thus, the K-condition number can be viewed as a useful tool for the construction of preconditionings. As soon as the preconditioning is specified, one can also try to estimate its standard (spectral) condition number in order to verify its efficiency. Several examples of such investigations are found in the paper.

## 3 Two-stage preconditionings

In this section, we will consider the following general scheme for preconditioning of SPD matrices. Let $A$ be a result of certain preprocessing of the original matrix $A_{0}$, e.g. by preordering, or scaling it to unit diagonal,

$$
A=\operatorname{Diag}\left(A_{0}\right)^{-1 / 2} A_{0} \operatorname{Diag}\left(A_{0}\right)^{-1 / 2}
$$

or, sometimes, even by a two-sided preconditioning by Incomplete Cholesky,

$$
A=U^{-T} A_{0} U^{-1}
$$

such as with the use of the IC2 preconditioning of [16].

- Stage 1. Let $Z$ be a nonsingular matrix with 2 by 2 block structure, and consider a congruence transformation of A, keeping the same block structure

$$
B=Z^{T} A Z=\left[\begin{array}{ll}
B_{11} & B_{12} \\
B_{21} & B_{22}
\end{array}\right]
$$

such that $K(B) \leq K(A)$. The main purpose of such a transformation is to reduce as much as possible the quantity

$$
\gamma=\left\|B_{11}^{-1 / 2} B_{12} B_{22}^{-1 / 2}\right\|
$$

which always satisfies $0<\gamma<1$. We will see that $\gamma$ should not be too close to 1 in order for the preconditioning to be efficient.

- Stage 2. Let $D$ be a block-diagonal 2 by 2 matrix with diagonal blocks $D_{1}$ and $D_{2}$ equal to the (approximate) inverses of $B_{11}$ and $B_{22}$, respectively. We shall refer to such a preconditioning as Approximate Block Jacobi preconditioning. Then, typically, $K(D B)<K(B) \leq K(A)$ and since $\lambda_{i}(D B)=\lambda_{1}(H A)$, it holds $K(D B)=K(H A)$, where the resulting preconditioner for $A$ will be

$$
H=Z D Z^{T}
$$

The effect of approximate Block Jacobi preconditioning on the spectal condition number and on the $K$-condition number were first studied in [4] and [15], respectively.

We consider first two illuminating examples of congruence transformations which can be related to the first stage of such preconditionings.

Example 1. Let

$$
Z=\left[\begin{array}{cc}
I_{1} & -A_{11}^{-1} A_{12} \\
0 & I_{2}
\end{array}\right]
$$

where $I_{i}, \quad i=1,2$ denote the identity matrices of consistent orders. Then an elementary computation shows that

$$
Z^{T} A Z=\left[\begin{array}{cc}
A_{11} & 0 \\
0 & S
\end{array}\right]
$$

where $S=A_{22}-A_{21} A_{11}^{-1} A_{12}$ is the Schur complement. Hence, in this case, $\gamma=0$. However, this is not a viable choice as it requires exact solutions of systems with $A_{11}$ and $S$, and $S$ is in general a full matrix.

It is anyhow interesting to note that the above transformation, when applied to the indefinite matrix

$$
A=\left[\begin{array}{cc}
A_{11} & A_{12} \\
A_{21} & 0
\end{array}\right]
$$

leads to

$$
Z^{T} A Z=\left[\begin{array}{cc}
A_{11} & 0 \\
0 & -A_{21} A_{11}^{-1} A_{12}
\end{array}\right]
$$

In certain applications (such as for certain finite element approximation of finite element Stokes problem for incompressible fluids), it turns out that the matrix $A_{21} A_{11}^{-1} A_{12}$ is well-conditioned and the transformation can hence be of interest in this case.

Example 2. We recall now (see e.g. [9] and [1]) another well known example of a congruence transformation showing the relation between the standard and hierarchical (nodal) basis function matrices. Let $J_{12}$ be the interpolation matrix between the sets of standard finite element basis functions and hierarchical basis functions, i.e., it holds $v_{S B}=J_{12} v_{H B}$ for corresponding elements in the two sets. The matrix $J_{12}$ is typically very sparse. Let $n_{2}$ be the number of degrees of freedom of the coarse space, let $n_{1}$ be that of the added basis functions, and let

$$
Z=\left[\begin{array}{cc}
I_{1} & J_{12} \\
0 & I_{2}
\end{array}\right]
$$

Thus if $A$ is the standard basis function matrix, it holds

$$
B=Z^{T} A Z=\left[\begin{array}{ll}
B_{11} & B_{12} \\
B_{21} & B_{22}
\end{array}\right]
$$

where

$$
B_{11}=A_{11}, \quad B_{12}=A_{12}+A_{11} J_{12}, \quad B_{21}=B_{12}^{T}
$$

and

$$
B_{22}=A_{22}+A_{21} J_{12}+J_{12}^{T} A_{12}+J_{12}^{T} A_{11} J_{12}
$$

Here $B$ is the hierarchical basis function matrix, see, e.g.[9]. While (e.g. in the case when a discretization of a 2-nd order elliptic equation is considered)

$$
1-\left\|A_{11}^{-1 / 2} A_{12} A_{22}^{-1 / 2}\right\|=O\left(h^{2}\right)
$$

where $h$ is a meshsize parameter, it holds that

$$
\gamma=\left\|B_{11}^{-1 / 2} B_{12} B_{22}^{-1 / 2}\right\|=1-c, \quad 0<c<1
$$

for some $c$ which does not depend on size, nor shape of elements and also not on jumps of coefficients if they occur only at the coarse mesh edges, for further details, see $[4,9]$.

While the hierarchical basis function submatrices $B_{12}, B_{21}$, are less sparse than the corresponding matrices for the standard basis function matrix $(A)$, the congruence transformation $Z^{T} A Z$ allows one to work with $A$ in computing actions of the iteration matrix.

Next we consider certain examples of Block Jacobi preconditionings, quite similar to those which were already described, e.g,, in [4], [1], [3].

### 3.1 Estimating the K-condition number for the Exact, Full and Partial <br> 2 by 2 Block Jacobi preconditioning

Let us first consider the simplest case of the exact Block Jacobi method with preconditioner

$$
H=\left[\begin{array}{cc}
B_{11}^{-1} & 0 \\
0 & B_{22}^{-1}
\end{array}\right]
$$

to the matrix $B$. It is well known that this preconditioning (up to an arbitrary positive scalar factor) is optimum over all 2 by 2 preconditionings with respect to the spectral condition number. Furthermore, the condition number of $H^{-1} B$ is $\kappa(H B)=$ $(1+\gamma) /(1-\gamma)$, where $\gamma=\left\|B_{11}^{-\frac{1}{2}} B_{12} B_{22}^{-\frac{1}{2}}\right\|$. The following result [15] (see also [1]) shows that such optimality holds also in the sense of the K-conditioning.

Theorem 3.1 Let an SPD $n \times n$ matrix $B$ be split into $2 \times 2$ block form as above and let $D_{1}$ and $D_{2}$ be arbitrary SPD matrices of the same orders $n_{1}$ and $n_{2}$ as $B_{11}$ and $B_{22}$, respectively, so that the block diagonal matrix

$$
D=\left[\begin{array}{cc}
D_{1} & 0 \\
0 & D_{2}
\end{array}\right]
$$

is also SPD. Then the minimum of the matrix functional $K(D B)$ is attained at $D_{1}=$ $B_{11}^{-1}$ and $D_{2}=B_{22}^{-1}$ and is equal to

$$
\min _{D_{1}, D_{2} \text { are } \mathrm{SPD}} K(D B)=K\left(D_{B}^{-1} B\right)=\frac{\operatorname{det}\left(B_{11}\right) \operatorname{det}\left(B_{22}\right)}{\operatorname{det}(B)}
$$

where

$$
D_{B}=\left[\begin{array}{cc}
B_{11} & 0 \\
0 & B_{22}
\end{array}\right]
$$

is the block diagonal part of $B$.

Proof. (See Section A1 of [15].) Since trace $(D B)=\operatorname{trace}\left(D D_{B}\right), n^{-1} \operatorname{trace}\left(D_{B}^{-1} B\right)=$ 1 and $\operatorname{det}(D B)=\operatorname{det}\left(D D_{B}\right) \operatorname{det}\left(D_{B}^{-1} B\right)$ it follows that the identity

$$
K(D B)=K\left(D D_{B}\right) K\left(D_{B}^{-1} B\right)
$$

holds. As follows from the arithmetic-geometric mean inequality, cf. [1, 15], the minimum of $K\left(D D_{B}\right)$ is equal to 1 and is attained if and only if $D D_{B}=\alpha I$ for some $\alpha>0$. Hence, $D=\alpha D_{B}^{-1}$ and the required result readily follows for $\alpha=1$. Q.E.D.

Remark 3.1 Clearly, the 2 by 2 splitting should be chosen such that $\operatorname{det}\left(B_{11} B_{22}\right)$ is as small as possible to obtain a better K-conditioning. Further, one can readily see that the attained value of $K$ is

$$
K\left(D_{B}^{-1} B\right)=1 / \operatorname{det}\left(I_{2}-C^{T} C\right), \quad C=B_{11}^{-1 / 2} B_{12} B_{22}^{-1 / 2}
$$

where $I_{2}$ is the identity matrix of order $n_{2}$, which stresses again the importance of making the norm of the matrix $C$ as small as possible.

In practice, it appears to be an important case when the matrix $D_{1}$ is prescribed, and only $D_{2}$ can be optimized. The corresponding preconditioning can be regarded as the Partial Block Jacobi one. The following result, which generalize that of Theorem 3.1, holds in this case.

Theorem 3.2 Let an SPD $n \times n$ matrix $B$ be split into $2 \times 2$ block form

$$
B=\left[\begin{array}{ll}
B_{11} & B_{12} \\
B_{21} & B_{22}
\end{array}\right]
$$

with the orders of the diagonal blocks $B_{11}$ and $B_{22}$ being $n_{1}$ and $n_{2}$, respectively. Let $D$ be the block diagonal matrix

$$
D=\left[\begin{array}{cc}
D_{1} & 0 \\
0 & D_{2}
\end{array}\right]
$$

with the $n_{1} \times n_{1} S P D$ block $D_{1}$ fixed and the $D_{2}$ block being an arbitrary $n_{2} \times n_{2} S P D$ matrix. Then the minimum of the matrix functional $K(D B)$ with respect to $D_{2}$ is attained for

$$
D_{2}=\sigma B_{22}^{-1}
$$

and is equal to

$$
\min _{D_{2} \text { is SPD }} K(D B)=\frac{\sigma^{n_{1}} \operatorname{det}\left(B_{22}\right)}{\operatorname{det}(B) \operatorname{det}\left(D_{1}\right)}=K\left(D_{1} B_{11}\right) K\left(D_{B}^{-1} B\right),
$$

where

$$
\sigma=\frac{1}{n_{1}} \operatorname{trace}\left(D_{1} B_{11}\right) .
$$

Proof. Using the inequality $K(X) \geq 1$ with $X=D_{2} B_{22}$ (which holds for any diagonalizable matrix $X$ with positive eigenvalues), one has

$$
\operatorname{trace}\left(D_{2} B_{22}\right) \geq n_{2}\left(\operatorname{det}\left(D_{2} B_{22}\right)\right)^{1 / n_{2}}
$$

Therefore, we obtain the following lower bound for $K(D B)$ :

$$
\begin{aligned}
& K(D B)=\frac{\left(\frac{1}{n}\left(\operatorname{trace}\left(D_{1} B_{11}\right)+\operatorname{trace}\left(D_{2} B_{22}\right)\right)\right)^{n}}{\operatorname{det}\left(D_{1}\right) \operatorname{det}\left(D_{2}\right) \operatorname{det}(B)} \geq \frac{\left(\frac{1}{n}\left(n_{1} \sigma+n_{2}\left(\operatorname{det}\left(D_{2} B_{22}\right)\right)^{1 / n_{2}}\right)\right)^{n}}{\operatorname{det}\left(D_{1}\right) \operatorname{det}\left(D_{2}\right) \operatorname{det}(B)} \\
& =\left(\frac{\frac{n_{1}}{n} \sigma+\frac{n_{2}}{n}\left(\operatorname{det}\left(D_{2} B_{22}\right)\right)^{1 / n_{2}}}{\left(\operatorname{det}\left(D_{2} B_{22}\right)\right)^{1 / n}}\right)^{n} \frac{\operatorname{det}\left(B_{22}\right)}{\operatorname{det}(B) \operatorname{det}\left(D_{1}\right)} \geq\left(\min _{\tau>0} \varphi(\tau)\right)^{n} \frac{\operatorname{det}\left(B_{22}\right)}{\operatorname{det}(B) \operatorname{det}\left(D_{1}\right)}
\end{aligned}
$$

where we denoted $\tau=\operatorname{det}\left(D_{2} B_{22}\right)$ and

$$
\varphi(\tau)=\frac{n_{1}}{n} \sigma \tau^{-\frac{1}{n}}+\frac{n_{2}}{n} \tau^{\frac{1}{n_{2}}-\frac{1}{n}} .
$$

An elementary computation shows now that the minimum of $\varphi$ equals $\sigma^{n_{1} / n}$ and is attained for

$$
\tau=\sigma^{n_{2}}
$$

Hence the expression for the optimum value of the K-condition number is proved. Further, it follows from the proof that this lower bound is attained when $K(X)=$ $K\left(D_{2} B_{22}\right)=1$, which yields $D_{2} B_{22}=\alpha I_{2}$ with $\alpha>0$. The above formula for $\tau$ now readily yields the required result by letting $\alpha=\sigma$. Q.E.D.

Remark 3.2 Theorem 3.2 shows that by first minimizing the K-condition number with respect to $D_{2}$ (for $D_{1}$ fixed) and then minimizing the resulting condition number with respect to $D_{1}$ the same optimality result holds as when the condition number is minimized by simultaneously varying $D_{1}$ and $D_{2}$.

Remark 3.3 In the case when both $D_{1}$ and $D_{2}$ are only approximations to the inverses of the diagonal blocks of $B$ but the following scaling property holds,

$$
\operatorname{trace}\left(D_{1} B_{11}\right) / n_{1}=\operatorname{trace}\left(D_{2} B_{22}\right) / n_{2}=1,
$$

a simple exact formula holds for the resulting K-condition number of the preconditioned matrix:

$$
\begin{equation*}
K(D B)=K\left(D_{1} B_{11}\right) K\left(D_{2} B_{22}\right) K\left(D_{B}^{-1} B\right) \tag{3.1}
\end{equation*}
$$

### 3.2 Estimating the spectral condition number for Approximate 2 by 2 Block Jacobi preconditioning

Let us now consider the estimates for the standard (spectral) condition number $\kappa(D A)$ obtained when applying an Approximate Block Jacobi preconditioning.

These results are found in $[4,1]$ but are presented here for completeness. In particular, we follow [1], pp. 378-380. They hold also for singular (i.e. positive semidefinite) matrices with $B_{11}$ nonsingular, for which $B v=0$ implies $B_{22} v_{2}=0$.

We consider then the extreme eigenvalues of the generalized eigenvalue problem $\lambda D x=B x$. Note that in this subsection, we let $D=\left[\begin{array}{cc}D_{1} & 0 \\ 0 & D_{2}\end{array}\right]$ where $D_{2}$ will be singular if $B_{22}$ is singular. Further, $\gamma$ is the constant in the strengthened Cauchy-Bunyakowski-Schwarz (CBS) inequality,

$$
x_{1}^{T} B_{12} x_{2} \leq \gamma\left\{x_{1}^{T} B_{11} x_{1} x_{2}^{T} B_{22} x_{2}\right\}^{\frac{1}{2}}
$$

If both $B_{11}$ and $B_{22}$ are positive definite, then, as we have seen, $\gamma=\left\|B_{11}^{-\frac{1}{2}} B_{12} B_{22}^{-\frac{1}{2}}\right\|$. Clearly, $\gamma<1$.
Below, the notation $A \geq B$ means that $A-B$ is positive semidefinite.

Theorem 3.3 Let $B$ be symmetric and positive semidefinite and split in a two by two block form such that if $B v=0$ then $B_{22} v_{2}=0$ when $v=\left[\begin{array}{l}v_{1} \\ v_{2}\end{array}\right]$ is split correspondingly. Let $\gamma$ be the constant in the corresponding strengthened CBS inequality and assume that

$$
\begin{aligned}
& \alpha_{1} B_{11} \leq D_{1} \leq \beta_{1} B_{11} \\
& \alpha_{2} B_{22} \leq D_{2} \leq \beta_{2} B_{22}
\end{aligned}
$$

for some $0<\alpha_{1} \leq \beta_{1}, 0<\alpha_{2} \leq \beta_{2}$. Then with $D=\left[\begin{array}{cc}D_{1} & 0 \\ 0 & D_{2}\end{array}\right]$,
a) $\lambda_{\max } \leq \frac{1}{\alpha_{1}}\left\{\frac{1}{2}\left(1+\frac{\alpha_{1}}{\alpha_{2}}\right)+\left[\left(\frac{1}{2}\left(1-\frac{\alpha_{1}}{\alpha_{2}}\right)\right)^{2}+\frac{\alpha_{1}}{\alpha_{2}} \gamma^{2}\right]^{\frac{1}{2}}\right\}$,

$$
\lambda_{\min } \geq \frac{1-\gamma^{2}}{\beta_{2}}\left\{\frac{1}{2}\left(1+\frac{\beta_{1}}{\beta_{2}}\right)+\left[\left(\frac{1}{2}\left(1-\frac{\beta_{1}}{\alpha_{2}}\right)\right)^{2}+\frac{\beta_{1}}{\beta_{2}} \gamma^{2}\right]^{\frac{1}{2}}\right\}^{-1}
$$

and the condition number of the preconditioned matrix $D^{-1} B$ is $\kappa \leq \lambda_{\max } / \lambda_{\min }$.
b) If we scale the blocks so that $\frac{\alpha_{1}}{\alpha_{2}} \leq \frac{\beta_{1}}{\beta_{2}}=1$, then $\kappa \leq \frac{\beta_{1}}{\beta_{2}} \frac{1+\gamma}{1-\gamma}$.
c) The following simplified upper bound holds,

$$
\kappa \leq \frac{1}{1-\gamma^{2}}\left(\frac{1}{\alpha_{1}}+\frac{1}{\alpha_{2}}\right)\left(\beta_{1}+\beta_{2}\right)
$$

Proof. The extreme eigenvalues are the extreme values of

$$
\frac{x^{T} B x}{x^{T} D x}=\frac{x_{1}^{T} B_{11} x_{1}+2 x_{1}^{T} B_{12} x_{2}+x_{2}^{T} B_{22} x_{2}}{x_{1}^{T} D_{1} x_{1}+x_{2}^{T} D_{2} x_{2}}
$$

Using the strengthened CBS-inequality and the arithmetic-geometric inequality $\sqrt{a b} \leq$ $\frac{1}{2}\left(\zeta a+\zeta^{-1} b\right)$, where $\zeta>0$, we find

$$
2\left|x_{1}^{T} B_{12} x_{2}\right| \leq \gamma \zeta x_{1}^{T} B_{11} x_{1}+\gamma \zeta^{-1} x_{2}^{T} B_{22} x_{2}
$$

This shows that

$$
\lambda_{\max } \leq \min _{\zeta>0} \max _{x_{1}, x_{2}} \frac{(1+\gamma \zeta) x_{1}^{T} B_{11} x_{1}+\left(1+\gamma \zeta^{-1}\right) x_{2}^{T} B_{22} x_{2}}{x_{1}^{T} D_{1} x_{1}+x_{2}^{T} D_{2} x_{2}}
$$

and using the given spectral relations between $B_{11}$ and $D_{1}$ and $B_{22}$ and $D_{2}$ we obtain

$$
\lambda_{\max } \leq \min _{\zeta>0} \max \left\{\frac{1+\gamma \zeta}{\alpha_{1}}, \frac{1+\gamma \zeta^{-1}}{\alpha_{2}}\right\}
$$

where the optimal value of $\zeta$ is found from the equation

$$
(1+\gamma \zeta) / \alpha_{1}=\left(1+\gamma \zeta^{-1}\right) / \alpha_{2}
$$

The lower eigenvalue bound is found in a similar way. Here

$$
\begin{aligned}
\lambda_{\min } & \geq \max _{\gamma<\zeta<\gamma-1} \min _{x_{1}, x_{2}} \frac{(1-\gamma \zeta) x_{1}^{T} B_{11} x_{1}+\left(1-\gamma \zeta^{-1}\right) x_{2}^{T} B_{22} x_{2}}{x_{1}^{T} D_{1} x_{1}+x_{2}^{T} D_{2} x_{2}} \\
& \geq \max _{\gamma<\zeta<\gamma^{-1}} \min \left\{\frac{1-\gamma \zeta}{\beta_{1}}, \frac{1-\gamma \zeta^{-1}}{\beta_{2}}\right\}
\end{aligned}
$$

where (if $\beta_{1} \leq \beta_{2}$, otherwise exchange $\zeta$ with $\zeta^{-1}$ above) the optimal value of $\zeta$ satisfies

$$
\gamma \zeta=\frac{1}{2}\left(1-\frac{\beta_{1}}{\beta_{2}}\right)+\left[\left(\frac{1}{2}\left(1-\frac{\beta_{1}}{\beta_{2}}\right)\right)^{2}+\frac{\beta_{1}}{\beta_{2}} \gamma^{2}\right]^{\frac{1}{2}}
$$

and we note that $\gamma^{2}<\gamma \zeta<1$, i.e. $\gamma<\zeta<\gamma^{-1}$, which gives the lower bound of $\lambda_{\min }$. Part b) follows by direct computation and to prove part c) we let $\gamma=1$ in the square root expressions for $\lambda_{\max }$ and $\lambda_{\min }$. Q.E.D.

### 3.3 Deriving the $n_{2}$-Rank Modification K-optimum preconditioning

Let us now consider preconditioners of the form

$$
H=I+V S V^{T}
$$

for the matrix $A$, where the matrix $V$ is a fixed $n \times n_{2}$ matrix with $n_{2} \ll n$, and $S$ is a symmetric $n_{2} \times n_{2}$ matrix depending properly on $V$ and $A$. This construction was investigated, e.g. in [12, 7, 18].

Following [12], let us choose $S$ as the solution of the following optimization problem:

$$
\begin{equation*}
S=\arg \min _{S=S^{T}} K\left(\left(I+V S V^{T}\right) A\right) \tag{3.2}
\end{equation*}
$$

The resulting preconditioning was referred to as Low Rank Modification (LRM) in [12], in view of the requirement of limiting of the size of $S$ when it is supposed to be calculated explicitly.
Assume that the columns of $V$ are orthogonalized and let

$$
Z_{2}=V\left(V^{T} V\right)^{-1 / 2}
$$

so that

$$
Z_{2}^{T} Z_{2}=I_{2}
$$

and introduce the $n \times n_{1}$ matrix $Z_{1}$ such that

$$
Z_{1}^{T} Z_{2}=0
$$

and

$$
Z_{1}^{T} Z_{1}=I_{1}
$$

In this case, the matrix

$$
Z=\left[Z_{1} Z_{2}\right]
$$

will be orthogonal, and, by $Z Z^{T}=I$, one has

$$
Z_{1} Z_{1}^{T}+Z_{2} Z_{2}^{T}=I
$$

One has then

$$
\begin{gathered}
K\left(\left(I+V S V^{T}\right) A\right)=K\left(Z^{T}\left(I+V S V^{T}\right) Z Z^{T} A Z\right)=K\left(\left(I+Z^{T} V S V^{T} Z\right)\left(Z^{T} A Z\right)\right) \\
=K\left(\left(I+Z^{T} Z_{2}\left(V^{T} V\right)^{1 / 2} S\left(V^{T} V\right)^{1 / 2} Z_{2}^{T} Z\right)\left(Z^{T} A Z\right)\right) \\
=K\left(\left[\begin{array}{cc}
I_{1} & 0 \\
0 & I_{2}+\left(V^{T} V\right)^{1 / 2} S\left(V^{T} V\right)^{1 / 2}
\end{array}\right]\left[\begin{array}{cc}
Z_{1}^{T} A Z_{1} & Z_{1}^{T} A Z_{2} \\
Z_{1}^{T} A Z_{2} & Z_{2}^{T} A Z_{2}
\end{array}\right]\right)
\end{gathered}
$$

Thus we have the same problem, the solution of which was given by Theorem 3.2. Therefore, setting

$$
\begin{gathered}
B_{i, j}=Z_{i}^{T} A Z_{j}, \quad i, j=1,2, \quad \sigma=\operatorname{trace}\left(Z_{1}^{T} A Z_{1}\right) / n_{1} \\
D_{1}=I_{1}, \quad D_{2}=I_{2}+\left(V^{T} V\right)^{1 / 2} S\left(V^{T} V\right)^{1 / 2}
\end{gathered}
$$

one has

$$
I_{2}+\left(V^{T} V\right)^{1 / 2} S\left(V^{T} V\right)^{1 / 2}=\sigma\left(Z_{2}^{T} A Z_{2}\right)^{-1}
$$

which gives

$$
\begin{gathered}
S=-\left(V^{T} V\right)^{-1}+\sigma\left(V^{T} V\right)^{-1 / 2}\left(Z_{2}^{T} A Z_{2}\right)^{-1}\left(V^{T} V\right)^{-1 / 2} \\
=-\left(V^{T} V\right)^{-1}+\sigma\left(\left(V^{T} V\right)^{1 / 2} Z_{2}^{T} A Z_{2}\left(V^{T} V\right)^{1 / 2}\right)^{-1} \\
=-\left(V^{T} V\right)^{-1}+\sigma\left(V^{T} A V\right)^{-1}
\end{gathered}
$$

where

$$
\sigma=\frac{\operatorname{trace}(A)-\operatorname{trace}\left(\left(V^{T} V\right)^{-1} V^{T} A V\right)}{n-n_{2}}
$$

This is the same formula as obtained in [12]:

$$
\begin{equation*}
H=\left(I-V\left(V^{T} V\right)^{-1} V^{T}\right)+\sigma V\left(V^{T} A V\right)^{-1} V^{T} \tag{3.3}
\end{equation*}
$$

In $[7,18]$ a somewhat different formula for $S$ (namely, without the term $-\left(V^{T} V\right)^{-1}$ and with a different choice of $\sigma$ ) was used.

Thereby, the following more general preconditioner was considered,

$$
\begin{equation*}
H=M^{-1}+\sigma V \widehat{C}^{-1} V^{T} \tag{3.4}
\end{equation*}
$$

where $M$ and $\widehat{C}$ are positive definite preconditioners for $A$ and $A_{V}=V^{T} A V$, respectively. Here $M$ is typically a smoother used to damp the higher eigenvalue modes of $A$ while $\widehat{C}$ can be chosen as a much simpler operator than $B_{V}$. The positive parameter $\sigma$ is chosen to move the set of smallest eigenvalues of $M^{-1} A$ to a cluster of bigger eigenvalues, in this way improving the conditioning significantly for ill-conditioned
problems where, typically, there exist several small eigenvalues of $A$. A good choice is $\sigma=\lambda_{\max }\left(M^{-1} A\right) / \lambda_{\max }\left(\widehat{C}^{-1} A_{V}\right)$, which number can normally be estimated with little expense (see [18]).

As we have seen, the projection operator from (3.3) is based on choosing $S$ to minimize the $K$-condition number in (3.2). However, as follows from the discussion in Section 2, this may not be the best choice in actually minimizing the number of iterations.

The following estimate of the extreme eigenvalues of $H A$ holds showing a significant reduction in the condition number when the vector space spanned by the column vectors of $V$ is sufficiently close to the eigenvector space for the smallest eigenvalues.

Theorem 3.4 ([18]) Let $H=M^{-1}+\sigma V \widehat{C}^{-1} V^{T}$ and assume that $\left\{\lambda_{i}, \underline{v}_{i}\right\}_{i=1}^{n}$ is an ordered set of eigenpairs of $M^{-1} A$ such that $\lambda_{1} \leq \ldots \leq \lambda_{n}$. Let the matrix $V_{e}=$ $\left[\underline{v}_{1}, \ldots, \underline{v}_{m}\right]$. If $V$ is such that the subspace $\mathcal{W}=\left(\operatorname{Im} A^{\frac{1}{2}} V\right)^{\perp}$ and $\mathcal{V}_{e}=\operatorname{Im}\left(A^{\frac{1}{2}} V_{e}\right)$ satisfy

$$
\gamma=\cos \left(\mathcal{W}, \mathcal{V}_{e}\right)=\sup _{\substack{x \in W \\ y \in \mathcal{V}_{e}}} \frac{x^{T} y}{\left\{x^{T} x y^{T} y\right\}^{\frac{1}{2}}}
$$

then the minimal eigenvalue of $C B$ is bounded as

$$
\lambda_{\min }(H A) \geq \max \left\{\lambda_{1},(1-\gamma) \min \left\{\frac{\lambda_{\max }\left(M^{-1} A\right)}{\kappa\left(\widehat{C}^{-1} A_{V}\right)}, \lambda_{m+1}\right\}\right\}
$$

and the maximal eigenvalue of $H A$ is bounded as

$$
\lambda_{\max }(H A) \leq 2 \lambda_{\max }\left(M^{-1} A\right)
$$

for any choce of $V$ and $\widehat{C}$.
As shown in [18], the theorem can be generalized to include eigenspaces of $\widehat{M}^{-1} \widehat{A}$ for nearby matrices $\widehat{M}$ and $\widehat{A}$ satisfying $\widehat{M} \geq M$ and $\widehat{A} \leq A$.

The preconditioning method (3.4) has been called approximate subspace projection (ASP) method. Note that in the next section we will consider a similar preconditioning with $M$ having rank $n_{1}$ and therefore presenting a generalization of the ASP and LRM preconditionings.

The numerical experiments presented in [12] showed that even with a relatively weak explicit preconditioning IIC the LRM techniques provides essential reduction of the total arithmetic costs of the method. (However, being used alone, LRM may even sometimes lead to a somewhat slower convergence.) Therefore, one may expect even greater improvements in the case of IC2-LRM preconditioning. Similarly, the extensive numerical experiments in [18] showed how the ASP method can be implemented in practice and gave several examples of significant reductions of the condition numers of several orders of magnitude.

As was pointed out in $[7,18]$ the subspace spanned by the columns of $V$ should well approximate the subspace corresponding to the eigenvectors of $A$ with the smallest eigenvalues. When there are few very small isolated eigenvalues of $A$, then such
a matrix $V$ can be computed via the Lanczos method, and the ASP preconditioning will give a substantial reduction of the iteration number even with small $n_{2}$. However, such matrices $V$ are not easily found in a general case, especially when $n_{2}$ may not be small. Some techniques are demonstrated in [18] to find a proper $V$ for discretizations of second order elliptic problems. In more general cases we will consider a somewhat different approach which can be related to approximate Schur complement type preconditioners.

Remark 3.4 It is interesting to note that the theory developed above in Section 3.2 can be applied here in order to estimate the spectral condition number of the LRM preconditioned matrix as defined above. One can readily see that in this case all the conditions of Theorem 3.3 hold with

$$
\xi_{1}=\lambda_{\min }\left(Z_{1}^{T} A Z_{1}\right), \quad \eta_{1}=\lambda_{\max }\left(Z_{1}^{T} A Z_{1}\right), \quad \xi_{2}=\eta_{2}=\sigma
$$

and

$$
\gamma=\left\|\left(Z_{1}^{T} A Z_{1}\right)^{-1 / 2} Z_{1}^{T} A Z_{2}\left(Z_{2}^{T} A Z_{2}\right)^{-1 / 2}\right\|
$$

Thus, the condition number estimate appears to be expressed essentially in the same terms as those of $[7,18]$ for a similar preconditioning.

Remark 3.5 Replacing this "exact" expression of $S$ satisfying (3.2) with a certain approximation

$$
S=-\left(V^{T} V\right)^{-1}+B_{V}^{-1}
$$

one should require, by Theorem 3.3, that

$$
\xi_{2}\left(I_{2}+\left(V^{T} V\right)^{1 / 2} S\left(V^{T} V\right)^{1 / 2}\right)^{-1} \leq Z_{2}^{T} A Z_{2} \leq \eta_{2}\left(I_{2}+\left(V^{T} V\right)^{1 / 2} S\left(V^{T} V\right)^{1 / 2}\right)^{-1}
$$

which is eqiuvalent to

$$
\xi_{2} B_{V} \leq V^{T} A V \leq \eta_{2} B_{V}
$$

Hence, if the latter spectral bounds holds, Theorem 3.3 gives an estimate for the spectral condition number attained with the Approximate LRM preconditioning

$$
H=\left(I-V\left(V^{T} V\right)^{-1} V^{T}\right)+V B_{V}^{-1} V^{T}
$$

In particular, one can see that as soon as $\eta_{2} / \xi_{2} \leq \eta_{1} / \xi_{1}$, such an approximation to $V^{T} A V$ can be regarded as quite acceptable.

## 4 Approximate Schur Complement type preconditionings

In this section, we present a common framework for two-level preconditionings using approximate Schur complements.

Let us suppose that the SPD matrix $A$ is preordered in a proper way and consider its $2 \times 2$ splitting as

$$
A=\left[\begin{array}{ll}
A_{11} & A_{12} \\
A_{21} & A_{22}
\end{array}\right]
$$

Let the orders of the diagonal blocks $A_{11}$ and $A_{22}$ be $n_{1}$ and $n_{2}$, respectively. For practical reasons, we assume that $n_{1} \gg n_{2} \gg 1$ and that the matrix $A_{11}$ is considerably better conditioned than $A$.
It is a well known fact that the following exact formula for the inverse matrix holds:

$$
A^{-1}=\left[\begin{array}{cc}
A_{11}^{-1}+A_{11}^{-1} A_{12} S^{-1} A_{21} A_{11}^{-1} & -A_{11}^{-1} A_{12} S^{-1} \\
-S^{-1} A_{21} A_{11}^{-1} & S^{-1}
\end{array}\right]
$$

where

$$
S=A_{22}-A_{21} A_{11}^{-1} A_{12}
$$

is the corresponding Schur complement. It is clarifying for the presentation to note that the above formula can be rewritten as a $n_{2}$-rank modification of a $n_{1}-\mathrm{rank}$ symmetric nonnegative definite matrix:

$$
A^{-1}=\left[\begin{array}{cc}
A_{11}^{-1} & 0 \\
0 & 0
\end{array}\right]+\left[\begin{array}{c}
-A_{11}^{-1} A_{12} \\
I_{2}
\end{array}\right] S^{-1}\left[\begin{array}{ll}
-A_{21} A_{11}^{-1} & I_{2}
\end{array}\right]
$$

Both of the above formulas are readily obtained from the following simple block matrix $L^{T} D L$-factorization:

$$
A^{-1}=\left[\begin{array}{cc}
I_{1} & -A_{11}^{-1} A_{12} \\
0 & I_{2}
\end{array}\right]\left[\begin{array}{cc}
A_{11}^{-1} & 0 \\
0 & S^{-1}
\end{array}\right]\left[\begin{array}{cc}
I_{1} & 0 \\
-A_{21} A_{11}^{-1} & I_{2}
\end{array}\right]
$$

As above, we denote by $I_{1}$ and $I_{2}$ the identity matrix of the order $n_{1}$ and $n_{2}$, respectively.
Another useful relation is $\operatorname{det}(A)=\operatorname{det}\left(A_{11}\right) \operatorname{det}(S)$.
The statement of the problem is rather simple: let us replace the matrix $A_{11}^{-1}$ by certain approximate inverses, symmetric $D_{1}$, or even unsymmetric $H_{1}$, that is,

$$
A_{11} \approx D_{1}^{-1}, \quad A_{11} H_{1} \approx I_{1}
$$

and determine the matrix $D_{2}$ (to be used instead of $S^{-1}$ ) in order to obtain the preconditioner

$$
H=\left[\begin{array}{cc}
D_{1} & 0  \tag{4.1}\\
0 & 0
\end{array}\right]+\left[\begin{array}{c}
-H_{1} A_{12} \\
I_{2}
\end{array}\right] D_{2}\left[\begin{array}{cc}
-A_{21} H_{1}^{T} & I_{2}
\end{array}\right]
$$

which is as close to $A^{-1}$ as possible, e.g. in the sense of minimization of $\kappa(H A)$ or $K(H A)$.
Note that if $D_{2}$ is dense, and $D_{1}=U_{1}^{-1} U_{1}^{-T}$ as is the case when an incomplete Cholesky decomposition for $A_{11}$ is used, and $H_{1}$ is taken as a sparse approximate inverse, then only about of

$$
2 \mathrm{nz}\left(U_{1}\right)+2 \mathrm{nz}\left(H_{1}\right)+2 \mathrm{nz}\left(A_{12}\right)+n_{2}^{2}
$$

floating point operations (flops) are needed to multiply such a preconditioner $H$ by a vector.
As follows from Theorem 3.2, if the additional scaling condition trace $\left(D_{1} A_{11}\right)=n_{1}$ holds, then the solution to the above problem is given by

$$
D_{2}=\left(A_{22}-A_{21}\left(H_{1}+H_{1}^{T}-H_{1}^{T} A_{11} H_{1}\right) A_{12}\right)^{-1}
$$

in the sense that this matrix gives the minimum value of both the spectral and the K-condition numbers, cf.[1, 15].
This can be easily demonstrated if one considers

$$
Z=\left[\begin{array}{cc}
I_{1} & -H_{1} A_{12} \\
0 & I_{2}
\end{array}\right]
$$

and writes the preconditioner as previously,

$$
H=Z D Z^{T}
$$

One has then

$$
K(H A)=K\left(Z D Z^{T} A\right)=K\left(D Z^{T} A Z\right)=K(D B)
$$

with

$$
B=Z^{T} A Z=\left[\begin{array}{cc}
A_{11} & \left(I_{1}-A_{11} H_{1}\right) A_{12} \\
A_{21}\left(I_{1}-H_{1}^{T} A_{11}\right) & A_{22}-A_{21}\left(H_{1}+H_{1}^{T}-H_{1}^{T} A_{11} H_{1}\right) A_{12}
\end{array}\right]
$$

Note that if a sparse approximate inverse $H_{1}$ is used, then the block $B_{22}$ is also sparse (or at least its rows are easily computable) which gives a possibility to use a relatively large $n_{2}$ and apply an approximate inversion also for the block $B_{22}$, e.g., using the IC2 factorization. Another possibility is to recursively apply the method, e.g. as in $[9,3]$.

Remark 4.1 Note that the same preconditioning (but with $D_{1}=H_{1}$ ) was cited in [17], formulas (2.9), (2.10), (2.12); the references therein go back to [22, 20]. In [17] the above formula for $H_{2}$ was found too complicated to be implemented in a multilevel method.

A similar construction was also used and analyzed in [21] (again with $D_{1}=H_{1}$, cf. formulas (2.8)-(2.10) there) with a reference to [19].

Remark 4.2 The obtained preconditioning appears to be rather similar to the approximate subspace projection method, ASP, or generalized augmented matrix method, GAM, see, e.g. [18] and references therein. Indeed, let us denote the $n \times n_{2}$ block by

$$
V=\left[\begin{array}{c}
-H_{1} A_{12} \\
I_{2}
\end{array}\right]
$$

then one has

$$
H=\left[\begin{array}{cc}
D_{1} & 0 \\
0 & 0
\end{array}\right]+V\left(V^{T} A V\right)^{-1} V^{T}
$$

However, the first term in GAM preconditioning was chosen to be of full rank $n$ rather than $n_{1}$ in our case. Such a restriction may likely impair the resulting preconditioning quality.

### 4.1 Improving K-conditioning by the Approximate Schur Complement

Let us consider an approach to the construction of the matrix $H_{1}$ approximating $A_{11}$ by the minimization of $K\left(Z^{T} A Z\right)$. Since $\operatorname{det}(Z)=1$, such a setting is actually reduced to

$$
\min _{H_{1} \text { is sparse }} \operatorname{trace}\left(A_{21}\left(I_{1}-A_{11} H_{1}\right)^{T} A_{11}^{-1}\left(I_{1}-A_{11} H_{1}\right) A_{12}\right),
$$

or, even simpler,

$$
\min _{H_{1} \text { is sparse }} \operatorname{trace}\left(-2 A_{21} H_{1} A_{12}+A_{21} H_{1}^{T} A_{11}^{-1} H_{1} A_{12}\right)
$$

which obviously presents an unconstrained quadratic optimization problem. The latter appears to be rather (structurally) complicated for general sparsity patterns of $H_{1}$; hence, let us consider certain special cases. Incidently, nearly the same minimization problem and similar constructions for the matrix $H_{1}$ were considered in [10].

In the case of a diagonal matrix $H_{1}=\operatorname{Diag}(h)$ the above optimization problem appears to be rather easily solvable (at least, approximately). One can find that the vector $h$ representing the diagonal of the matrix $H_{1}$, can be found as the solution of the system

$$
\left(A_{11} \circ\left(A_{12} A_{21}\right)\right) h=\operatorname{diag}\left(A_{12} A_{21}\right)
$$

where " $\circ$ " stands for the Hadamard (componentwise) product of matrices. The matrix of this system typically has strong diagonal dominance, which makes it possible to determine an approximation to $h$ by a simple iterative method. It turns out that the case when $A_{12} A_{21}$ has zero diagonal entries yields virtually no essential complications.

In the case, when the matrix $H_{1}$ is chosen as a polynomial in $A_{11}$,

$$
H_{1}=q_{k-1}\left(A_{11}\right)=\sum_{i=1}^{k} \gamma_{i} A_{11}^{i-1}
$$

the polynomial coefficients can be found as the solution of the Hankel type system

$$
\left[\begin{array}{cccc}
\mu_{1} & \mu_{2} & \cdots & \mu_{k} \\
\mu_{2} & \mu_{3} & \cdots & \mu_{k+1} \\
\cdots & \cdots & \cdots & \cdots \\
\mu_{k} & \mu_{k+1} & \cdots & \mu_{2 k-1}
\end{array}\right]\left[\begin{array}{c}
\gamma_{1} \\
\gamma_{2} \\
\cdots \\
\gamma_{k}
\end{array}\right]=\left[\begin{array}{c}
\mu_{0} \\
\mu_{1} \\
\cdots \\
\mu_{k-1}
\end{array}\right]
$$

where

$$
\mu_{i}=\operatorname{trace}\left(A_{21} A_{11}^{i} A_{12}\right)
$$

Of course, $k$ should not be large in order to make the values of $\mu_{i}$ and the entries of $B_{22}$ easily computable.

An important property of this Approximate Schur Preconditioning is that it tends to improve the value of $K\left(D_{B}^{-1} B\right)$ as compared to $K\left(D_{A}^{-1} A\right)$, which seems important in view of the relation (3.1). Indeed, one has

$$
K\left(D_{B}^{-1} B\right)=1 / \operatorname{det}\left(D_{B}^{-1} B\right)=1 / \operatorname{det}\left(I_{1}-B_{11}^{-1 / 2} B_{12} B_{22}^{-1} B_{21} B_{11}^{-1 / 2}\right)
$$

$$
\begin{aligned}
& \text { where } \\
& \qquad B_{11}^{-1 / 2} B_{12} B_{22}^{-1} B_{21} B_{11}^{-1 / 2} \\
& =A_{11}^{-1 / 2}\left(I_{1}-A_{11} H_{1}\right) A_{12}\left(A_{22}-A_{21}\left(H_{1}+H_{1}^{T}-H_{1}^{T} A_{11} H_{1}\right) A_{12}\right)^{-1} A_{21}\left(I_{1}-H_{1}^{T} A_{11}\right) A_{11}^{-1 / 2} \\
& =A_{11}^{-1 / 2}\left(I_{1}-A_{11} H_{1}\right) A_{12}\left(S+A_{21}\left(I_{1}-H_{1}^{T} A_{11}\right) A_{11}^{-1}\left(I_{1}-A_{11} H_{1}\right) A_{12}\right)^{-1} A_{21}\left(I_{1}-H_{1}^{T} A_{11}\right) A_{11}^{-1 / 2}
\end{aligned}
$$

Then, using the equality

$$
\operatorname{det}\left(I_{1}-X\left(S+X^{T} X\right)^{-1} X^{T}\right)=\operatorname{det}(S) / \operatorname{det}\left(S+X^{T} X\right)
$$

with

$$
X=A_{11}^{-1 / 2}\left(I_{1}-A_{11} H_{1}\right) A_{12}
$$

one gets

$$
\begin{aligned}
& K\left(D_{B}^{-1} B\right)=\operatorname{det}\left(I_{2}+S^{-1 / 2} A_{21}\left(I_{1}-H_{1}^{T} A_{11}\right) A_{11}^{-1}\left(I_{1}-A_{11} H_{1}\right) A_{12} S^{-1 / 2}\right) \\
& \left.\quad \leq\left(1+\frac{\left\|S^{-1}\right\|}{n_{2}} \operatorname{trace}\left(A_{21}\left(I_{1}-H_{1}^{T} A_{11}\right) A_{11}^{-1}\left(I_{1}-A_{11} H_{1}\right) A_{12}\right)\right)\right)^{n_{2}}
\end{aligned}
$$

Note that the latter estimate actually presents an upper bound for $K\left(D_{B}^{-1} B\right)$ in terms of $K(B)$, the minimization of which with respect to $H_{1}$ has been discussed in this subsection.

### 4.2 Improving spectral conditioning by the Approximate Schur Complement

With respect to the estimation of the spectral condition number, one may expect a considerable reduction of $\gamma$ as compared to the original matrix $A$. For instance, it can be shown that if

$$
A_{21}\left(A_{11}^{-1}-H_{1}-H_{1}^{T}+H_{1}^{T} A_{11} H_{1}\right) A_{12} \leq \rho^{2} A_{21} A_{11}^{-1} A_{12}, \quad \rho<1
$$

then

$$
\frac{\gamma_{B}^{2}}{1-\gamma_{B}^{2}} \leq \rho^{2} \frac{\gamma_{A}^{2}}{1-\gamma_{A}^{2}}
$$

where

$$
\gamma_{A}=\left\|A_{11}^{-1 / 2} A_{12} A_{22}^{-1 / 2}\right\|, \quad \gamma_{B}=\left\|B_{11}^{-1 / 2} B_{12} B_{22}^{-1 / 2}\right\| .
$$

The proof can easily be constructed using the same formulas as in the end of the preceeding subsection. In particular, one can see that

$$
\gamma_{B}^{2}=\max _{y \neq 0} \frac{y^{T} X^{T} X y}{y^{T} S y+y^{T} X^{T} X y}
$$

with the same $X$ as above, and therefore, by $X^{T} X \leq \rho^{2} A_{21} A_{11}^{-1} A_{12}$,

$$
\gamma_{B}^{2} \leq \frac{\omega}{1+\omega}, \quad \omega=\lambda_{\max }\left(S^{-1} A_{21} A_{11}^{-1} A_{12}\right)=\frac{\gamma_{A}^{2}}{1-\gamma_{A}^{2}}
$$

Hence, the required estimate readily follows.

### 4.3 The choice of 2 by 2 splitting of the coefficient matrix

As was demonstrated above, it is advantageous to have the block $A_{11}$ not only of large size $n_{1}$ but as well-conditioned as possible, since the latter requirement makes it easier to find a good approximate inverse for it. Also, it is advantageous when the columns of $A_{21}$ are pairwise orthogonal, or nearly orthogonal. The latter condition can easily be satisfied if $A$ is a sparse matrix, e.g. of the type arising when solving boundary value problems for elliptic PDE's using FD or FE discretizations. Hence, the splitting can be based on the extraction of the block $A_{22}$ corresponding to an "independent set" of grid nodes. Otherwise, when $A$ is not sparse or its sparsity is not regular enough, one can base the splitting of $A$ using a certain "threshold pivot" Incomplete Cholesky factorization. For instance, supposing that $A$ is symmetrically scaled to unit diagonal, one can set a certain threshold parameter $(\theta<1$ and, in the course of an incomplete factorization (e.g. IC2 algorithm [16]) at the $k$-th step, one sets the whole current column of the right Cholesky factor $U$ equal to the $k$-th column of the identity matrix whenever it appears that the actually computed value is $u_{i i} \leq \theta$. Such an algorithm returns the IC 2 factorization of certain submatrix $A_{11}$ of $A$ such that all diagonal elements of its IC factor are sufficiently close to 1 . The value of $n_{1}$ can be adjusted by a proper choice of the threshold $\theta$.

## 5 The Preconditioned CGNR method with two-sided 2-norm error bounds

A possible application of high quality preconditionings for general SPD matrices may be the iterative solution of highly unsymmetric (sparse) linear systems.

In order to find the solution of unsymmetric linear system

$$
A_{0} x=b
$$

let us use the substitution $x=A_{0}^{T} y$; the system then takes the form

$$
A y=b
$$

with

$$
A=A_{0} A_{0}^{T}
$$

Using now the PCG algorithm presented in [6] and using there the substitutions $x_{k}=A_{0} y_{k}, q_{k}=A_{0}^{T} p_{k}$, one readily gets the solution method in the following form:

$$
r_{0}=b-A_{0} x_{0}, \quad q_{0}=A_{0}^{T} H r_{0} ; \quad \sigma_{-1}=0,
$$

$$
\begin{aligned}
\text { for } i= & 0,1, \ldots: \\
\sigma_{i}= & \sigma_{i-1}+1 / r_{i}^{T} H r_{i}, \\
& \text { if }(i>0 \text { and } i=0 \bmod d) \text { then } \\
& \quad \text { if }\left(2\left\|r_{i}\right\| \leq\left\|b-A_{0} x_{i}\right\|\right) \text { then QUIT_1 } \\
& \mu_{1}^{(i)}=\lambda_{\min }\left(T_{i}\right) \\
& \quad \text { if }\left(\mu_{1}^{(i)} \approx \mu_{1} \text { and } \sum_{j=i-d}^{i-1} \omega_{j}+\frac{1}{\sigma_{i} \mu_{1}^{(i)}} \leq \frac{\varepsilon^{2}}{1-\varepsilon^{2}} \sum_{j=0}^{i-1} \omega_{j}\right) \text { then QUIT_2 } \\
& \quad \text { endif } \\
\alpha_{i}= & r_{i}^{T} H r_{i} / q_{i}^{T} q_{i}, \\
\omega_{i}= & \left(r_{i}^{T} H r_{i}\right)^{2} / q_{i}^{T} q_{i}, \\
x_{i+1}= & x_{i}+q_{i} \alpha_{i}, \\
r_{i+1}= & r_{i}-A_{0} q_{i} \alpha_{i}, \\
\beta_{i}= & r_{i+1}^{T} H r_{i+1} / r_{i}^{T} H r_{i}, \\
q_{i+1}= & H r_{i+1}+q_{i} \beta_{i},
\end{aligned}
$$

where

$$
T_{i}=\left[\begin{array}{cccc}
\frac{1}{\alpha_{0}} & -\frac{\beta_{0}}{\alpha_{0}} & 0 & \ldots \\
-\frac{1}{\alpha_{0}} & \frac{1}{\alpha_{1}}+\frac{\beta_{0}}{\alpha_{0}} & -\frac{\beta_{1}}{\alpha_{1}} & \cdots \\
\ldots \ldots & \ldots & \ldots & \ldots \\
\cdots & 0 & -\frac{1}{\alpha_{i-2}} & \frac{1}{\alpha_{i-1}}+\frac{1}{\alpha_{i-2}}
\end{array}\right]
$$

and it is assumed that the minimum eigenvalue of the latter tridiagonal matrix is sufficiently close to

$$
\mu_{1}=\lambda_{\min }(H A) .
$$

The latter situation typically takes place for sufficiently good preconditionings and small values of $\varepsilon$ as the CG iterations enter the final stage of superlinear convergence.

Using the results of [6] one can see, from

$$
\left\|y-y_{i}\right\|_{A}=\left\|x-x_{i}\right\|_{,}
$$

that the following two-sided estimate holds:

$$
\sum_{j=i-d}^{i-1} \omega_{j} \leq\left\|x-x_{i-d}\right\|^{2} \leq \sum_{j=i-d}^{i-1} \omega_{j}+\frac{\mu_{1}^{-1}}{\sum_{j=0}^{i} \gamma_{j}^{-1}}, \quad d \leq i \leq n .
$$

In practice, a value such as $d=10$ proved to be satisfactory. When returning by QUIT_2, one can see that the above upper bound, the equality

$$
\left\|x-x_{k}\right\|^{2}=\sum_{j=k}^{n-1} \omega_{j}, \quad k=0, \ldots, n-1,
$$

used with $k=0$ and $k=i-d$, and the well-known inequality $\left\|x-x_{i}\right\| \leq\left\|x-x_{i-d}\right\|$ guarantee that the iterations are terminated with

$$
\left\|x-x_{i}\right\| \leq \varepsilon\left\|x-x_{0}\right\| .
$$

The return by QUIT_1 corresponds to an inacceptably large discrepancy between the iterated residual and the "exact" residual. In this case, the restart of CGNE with $x_{0}:=x_{i}$ should be performed in an attempt to achieve the required precision.

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