DELFT UNIVERSITY OF TECHNOLOGY<br>Faculty of Electrical Engineering, Mathematics and Computer Science

## ANSWERS OF THE TEST SCIENTIFIC COMPUTING ( wi4201 ) <br> Friday February 1 2019, 13:30-16:30

This are short answers, which indicate how the exercises can be answered. In most of the cases more details are needed to give a sufficiently clear answer.

1. (a) No.

Note that $\|A\|_{2}=\max _{1 \leq i \leq n}\left|\lambda_{i}\right|$ where $\lambda_{i}$ is an eigenvalue of $A$. Counter example

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

Note that $\max _{1 \leq i \leq n} \lambda_{i}=-1$, whereas $\|A\|_{2}=2$.
(b) Yes

From SPD it follows that $\mathbf{x}^{T} A \mathbf{x}>0$. Take the following vectors: $\mathbf{x}=\mathbf{e}_{i}$, where $\mathbf{e}_{i}$ is the $i^{\text {th }}$ column of the identity matrix $I$. This implies that

$$
\mathbf{e}_{i}^{T} A \mathbf{e}_{i}=a_{i i}>0
$$

which proves the claim.
(c) No

From $\mathbf{r}=\alpha_{1} \mathbf{v}_{1}+\alpha_{2} \mathbf{v}_{2}$ it follows that $A^{k} \mathbf{r}=c_{1} \mathbf{v}_{1}+c_{2} \mathbf{v}_{2}$. So all powers of $A$ multiplied with $\mathbf{r}$ are element of the span of $\mathbf{v}_{1}, \mathbf{v}_{2}$. This implies that the dimension of $K^{10}(A, \mathbf{r})$ is at most equal to 2 .
(d) Yes

Note that the following inequality is valid:

$$
\|\mathbf{u}\|_{2}=\sqrt{\sum_{i=1}^{n}\left(u_{i}\right)^{2}} \leq \sqrt{n \max _{1 \leq i \leq n}\left|u_{i}\right|^{2}}=\sqrt{n}\|\mathbf{u}\|_{\infty}
$$

(e) No
$\|A\|_{\infty}$ is the maximal absolute row sum, and $\|A\|_{1}$ is the maximal absolute column sum. In general these are not equal. Counter example

$$
A=\left[\begin{array}{ll}
1 & 1 \\
0 & 0
\end{array}\right]
$$

$$
\|A\|_{\infty}=2,\|A\|_{1}=1
$$

2. (a) For the discretization we use the Finite Difference Method on the given grid. Using nearest neighbours we have for the internal nodes that

$$
\begin{equation*}
\frac{\partial^{2} u}{\partial x^{2}}\left(x_{i}, y_{j}\right)=\frac{u_{i-1, j}^{h}-2 u_{i, j}^{h}+u_{i+1, j}^{h}}{h^{2}} \mathcal{O}\left(h^{2}\right) \text { for } 2 \leq i, j \leq m \tag{1}
\end{equation*}
$$

(and similar for the $y$-derivative). Using Taylor polynomials the claim that the error is $\mathcal{O}\left(h^{2}\right)$ should be shown. The approximation to the partial differential equation discretized on internal points of the grid can be written as

$$
\begin{equation*}
\frac{-u_{i, j-1}^{h}-u_{i-1, j}^{h}+4 u_{i, j}^{h}-u_{i+1, j}^{h}-u_{i, j+1}^{h}}{h^{2}}+u_{i, j}^{h}=f_{i, j}^{h} \text { for } 2 \leq i, j \leq m . \tag{2}
\end{equation*}
$$

(b) The stencil on the internal nodes is given by:

$$
\frac{1}{h^{2}}\left[\begin{array}{ccc}
0 & -1 & 0  \tag{3}\\
-1 & 4+h^{2} & -1 \\
0 & -1 & 0
\end{array}\right]
$$

The stencil in the lower left corner is:

$$
\frac{1}{h^{2}}\left[\begin{array}{ccc}
0 & -1 & 0  \tag{4}\\
0 & 4+h^{2} & -1 \\
0 & 0 & 0
\end{array}\right]
$$

(c) The matrix $A$ has the following structure:

$$
\frac{1}{h^{2}}\left(\begin{array}{cccccc}
T^{h} & -I & 0 & \ldots & \ldots & 0 \\
-I & T^{h} & -I & 0 & \ldots & 0 \\
\ldots & \ldots & \ldots & \ldots & \ldots & \ldots \\
0 & \ldots & 0 & -I & T^{h} & -I \\
0 & \ldots & \ldots & 0 & -I & T^{h}
\end{array}\right)
$$

where $T^{h}$ is given by

$$
\frac{1}{h^{2}}\left(\begin{array}{cccccc}
4+h^{2} & -1 & 0 & \ldots & \ldots & 0 \\
-1 & 4+h^{2} & -1 & 0 & \ldots & 0 \\
\ldots & \ldots & \ldots & \ldots & \ldots & \ldots \\
0 & \ldots & 0 & -1 & 4+h^{2} & -1 \\
0 & \ldots & \ldots & 0 & -1 & 4+h^{2}
\end{array}\right)
$$

Finally, the bandwidth is equal to $m$.
(d) That the matrix is symmetric is easy to explain. This implies that all eigenvalues are real valued. Using Gershgorin's theorem for the matrix, implies that all eigenvalues are larger than 1 . This implies that the matrix is SPD, see the line below Gershgorin's theorem in the lecture notes. Gershgorin's theorem is given
by: (Gershgorin) If $\lambda \in \sigma(A)$, then $\lambda$ is located in one of the $n$ closed disks in the complex plane that has center $a_{i i}$ and radius

$$
\begin{equation*}
\rho_{i}=\sum_{j=1, j \neq i}^{n}\left|a_{i j}\right| \tag{5}
\end{equation*}
$$

i.e.,

$$
\begin{equation*}
\lambda \in \sigma(A) \Rightarrow \exists i \text { such that }\left|a_{i i}-\lambda\right| \leq \rho_{i} . \tag{6}
\end{equation*}
$$

(e) The best Krylov subspace method for SPD matrices is: Conjugate Gradients (CG). Motivation for this choice: CG has short recurrences, its approximation is an element of the Krylov subspace and is has a minimization property that the error is minimal measured in the $A$ norm.
3. (a) The linear system $A \mathbf{u}=\mathbf{f}$ can then be written as $M \mathbf{u}=N \mathbf{u}+\mathbf{f}$. By multiplying to the left and right by $M^{-1}$ we can define an iterative scheme

$$
\begin{align*}
\mathbf{u}^{k+1} & =M^{-1} N \mathbf{u}^{k}+M^{-1} \mathbf{f} \\
& =M^{-1}(M-A) \mathbf{u}^{k}+M^{-1} \mathbf{f} \\
& =\mathbf{u}^{k}+M^{-1}\left(\mathbf{f}-A \mathbf{u}^{k}\right) \\
& =\mathbf{u}^{k}+M^{-1} \mathbf{r}^{k} \tag{7}
\end{align*}
$$

The recursion for the residual vector is given by

$$
\begin{align*}
\mathbf{r}^{k+1} & =\mathbf{f}-A \mathbf{u}^{k+1}  \tag{8}\\
& =\mathbf{f}-A \mathbf{u}^{k}-A M^{-1} \mathbf{r}^{k} \\
& =\mathbf{r}^{k}-A M^{-1} \mathbf{r}^{k} \\
& =\left(I-A M^{-1}\right) \mathbf{r}^{k}
\end{align*}
$$

(b) The error $\mathbf{e}^{k}$ and the residual $\mathbf{r}^{k}$ are related in the following way $\mathbf{r}^{k}=A \mathbf{e}^{k}$. We can use this relation to define an iterative scheme by the following sequence of three steps

- compute the defect: $\mathbf{r}^{k}=\mathbf{f}-A \mathbf{u}^{k}$;
- compute the approximate correction by solving the approximate residual equations: $\widehat{A} \widehat{\mathbf{e}}^{k}=\mathbf{r}^{k}$;
- add the correction to the previous iterand $\mathbf{u}^{k+1}=\mathbf{u}^{k}+\widehat{\mathbf{e}}^{k}$.
(c) Note that the iteration matrix is given by: $B=I-M^{-1} A$. This implies that $B_{\text {RIC }}=I-\tau A$. The resulting stencil is:

$$
\left[B_{R I C}\right]=\left[\begin{array}{lll}
\frac{\tau}{h^{2}} & 1-2 \frac{\tau}{h^{2}} & \frac{\tau}{h^{2}}
\end{array}\right]
$$

The stencil for $M_{G S}$ is easy to find:

$$
\left[M_{G S}\right]=\left[\begin{array}{lll}
-\frac{1}{h^{2}} & \frac{2}{h^{2}} & 0
\end{array}\right]
$$

However the stencil for $B_{G S}$ is not easy to find.
(d) Using the definition of $B_{R I C}=I-\tau A$ it easily follows that the eigenvalues of $B_{R I C}$ are given by the following expression: $\lambda_{R I C, i}=1-\tau \lambda_{i}$ where $\lambda_{i}$ are the eigenvalues of $A$. For the maximal absolute eigenvalue of $B_{R I C}$ (the spectral radius) we note that $\max _{1 \leq i \leq n}\left|\lambda_{R I C, i}\right|=\max \left[\left(1-\tau \lambda_{1}\right),\left(\tau \lambda_{n}-1\right)\right]$. One obtains the best speed up if these values are equal, so $1-\tau \lambda_{1}=\tau \lambda_{n}-1$. Solving this leads to the optimal value: $\tau=\frac{2}{\lambda_{1}+\lambda_{n}}$
(e) Determine the spectral radius of the Richardson iteration matrix. Note that $1-\tau \lambda_{1}=\frac{19}{20}$ and $1-\tau \lambda_{n}=\frac{1}{2}$, so the spectral radius is $\frac{19}{20}$. So if $\left(\frac{19}{20}\right)^{k} \leq 10^{-4}$ where $k$ is the number of iterations we have satisfied the stopping criterion. So the answer is: $k=\frac{-4}{\log \left(\frac{19}{20}\right)}=180$
4. (a) We take $\mathbf{u}^{1}=\alpha_{0} \mathbf{r}^{0}$ where $\alpha_{0}$ is a constant which has to be chosen such that $\left\|\mathbf{f}-A \mathbf{u}^{1}\right\|_{A}$ is minimal. This leads to

$$
\left\|\mathbf{f}-A \mathbf{u}^{1}\right\|_{A}^{2}=\left(\mathbf{f}-\alpha_{0} A \mathbf{r}^{0}\right)^{T} A\left(\mathbf{f}-\alpha_{0} A \mathbf{r}^{0}\right)=\mathbf{f}^{T} A \mathbf{f}-2 \alpha_{0}\left(A \mathbf{r}^{0}\right)^{T} A \mathbf{f}+\alpha_{0}^{2}\left(A \mathbf{r}^{0}\right)^{T} A A \mathbf{r}_{0} .
$$

The norm given above is minimized if $\alpha_{0}=\frac{\left(A \mathbf{r}^{0}\right)^{T} A \mathbf{f}}{\left(A \mathbf{r}^{0}\right)^{T} A^{2} \mathbf{r}^{0}}$.
(b) The optimality property of CG implies that the approximation $\mathbf{u}^{k}$ coming from CG satisfies:

$$
\left\|\mathbf{u}-\mathbf{u}^{k}\right\|_{A}=\min _{\mathbf{y} \in K^{k}\left(A ; \mathbf{r}^{0}\right)}\|\mathbf{u}-\mathbf{y}\|_{A}
$$

If the method terminates before we reach $k=n$ we know that we have a 'lucky' breakdown so $\mathbf{u}^{k}=u$. If not we know that the dimension of $K^{n}$ is equal to $n$, thus $K^{n}=\mathbb{R}^{n}$ and thus $\mathbf{u}^{n}=u$.
(c) The convergence of CG depends on the condition number. For SPD matrices the condition number is defined as

$$
K_{2}(A)=\frac{\lambda_{n}}{\lambda_{1}} .
$$

For a smaller condition number the convergence of CG is faster. Since $K_{2}\left(A_{1}\right)=$ 10 and $K_{2}\left(A_{2}\right)=200$, it is clear that we expect that the convergence for $A_{1}$ is much faster than for $A_{2}$.
(d) The superlinear convergence is given in the figure below: The explanation for


Figure 1: A superlinear convergent behavior
superlinear convergence is that initially the convergence is determined by the condition number $K_{2}(A)=\frac{\lambda_{n}}{\lambda_{1}}$. However after a number of iteration the effect of the smallest eigenvalue component is no longer influencing the convergence, so the convergence is determined by the effective condition number: $\frac{\lambda_{n}}{\lambda_{2}}$. So the convergence becomes faster and faster.
(e) The three properties are:
i. The matrix $M$ should be SPD.
ii. the eigenvalues of $M^{-1} A$ should be clustered around 1 , or the condition number of $M^{-1} A$ is (much) smaller than the condition number of $A$.
iii. it should be possible to obtain $M^{-1} \mathbf{y}$ at a low cost.
5. (a) In order to solve the linear system $A \mathbf{u}=\mathbf{f}$ with LU-decomposition without pivotting, we do the following steps:

- Find a lower triangular matrix $L$ and an upper triangular matrix $U$, such that $L U=A$ and the diagonal elements of $L$ are equal to 1 .
- Solve $\mathbf{y}$ from $L \mathbf{y}=\mathbf{f}$.
- Solve $\mathbf{u}$ from $U \mathbf{u}=\mathbf{y}$.

Since $L$ and $U$ are triangular matrices, this solution process is easy to implement. For the derivation of the costs see the lecture notes. The answer for a full matrix is for the decomposition the cost is $\frac{2}{3} n^{3}$ and for both solution steps together $2 n^{2}$.
(b) If we do the multiplication:

$$
\left(I-\alpha^{(k)} \mathbf{e}_{k}^{T}\right)\left(I+\alpha^{(k)} \mathbf{e}_{k}^{T}\right)
$$

we obtain the following:

$$
\begin{gathered}
I-\alpha^{(k)} \mathbf{e}_{k}^{T}+\alpha^{(k)} \mathbf{e}_{k}^{T}+\alpha^{(k)} \mathbf{e}_{k}^{T} \alpha^{(k)} \mathbf{e}_{k}^{T}= \\
I+\alpha^{(k)} \mathbf{e}_{k}^{T} \alpha^{(k)} \mathbf{e}_{k}^{T}
\end{gathered}
$$

Due to the zero structure of $\mathbf{e}_{k}$ and $\alpha^{(k)}$ the product $\mathbf{e}_{k}^{T} \alpha^{(k)}$ is equal to zero, so the last term is equal to zero, so

$$
\left(I-\alpha^{(k)} \mathbf{e}_{k}^{T}\right)\left(I+\alpha^{(k)} \mathbf{e}_{k}^{T}\right)=I
$$

which proves the claim that $M_{k}^{-1}=I+\alpha^{(k)} \mathbf{e}_{k}^{T}$.
(c) The perturbed solution $\mathbf{u}+\Delta \mathbf{u}$ solves the system

$$
\begin{equation*}
A(\mathbf{u}+\Delta \mathbf{u})=\mathbf{f}+\Delta \mathbf{f} \tag{9}
\end{equation*}
$$

Due to linearity, the perturbation $\Delta \mathbf{u}$ then solves the system

$$
\begin{equation*}
A \Delta \mathbf{u}=\Delta \mathbf{f} \tag{10}
\end{equation*}
$$

from which $\Delta \mathbf{u}=A^{-1} \Delta \mathbf{f}$ and therefore $\|\Delta \mathbf{u}\| \leq\left\|A^{-1}\right\|\|\Delta \mathbf{f}\|$. It follows form the multiplicative property that $\|\mathbf{f}\| \leq\|A\|\|\mathbf{u}\|$ and therefore

$$
\begin{equation*}
\frac{1}{\|\mathbf{u}\|} \leq\|A\| \frac{1}{\|\mathbf{f}\|} \tag{11}
\end{equation*}
$$

Combining these inequalities we arrive at the following bound on the norm of the perturbed solution

$$
\begin{equation*}
\frac{\|\Delta \mathbf{u}\|}{\|\mathbf{u}\|} \leq\left\|A^{-1}\right\|\|A\| \frac{\|\Delta \mathbf{f}\|}{\|\mathbf{f}\|}=\kappa(A) \frac{\|\Delta \mathbf{f}\|}{\|\mathbf{f}\|} \leq \delta \kappa(A) \tag{12}
\end{equation*}
$$

where $\kappa(A)$ denotes the condition number of $A$ measured in the norm $\|\cdot\|$.
(d) From the construction of $L$ and $U$ it follows that there are only zeroes outside the band with bandwidth $m$. Within the band, elements which are zero in $A$ become in general non-zero in $L$ and $U$ due to fill in.
(e) The difference is that we now look for a matrix $G$ such that $A=G G^{T}$, where $G$ is a lower triangular system. Advantages are:
i. Only half of the memory is needed.
ii. Only half of the amount of flops is needed.
iii. The method is stable so no pivotting is needed.

