DELFT UNIVERSITY OF TECHNOLOGY FACULTY OF ELECTRICAL ENGINEERING, MATHEMATICS AND COMPUTER SCIENCE

ANSWERS OF THE TEST SCIENTIFIC COMPUTING (wi4201) Wednesday January 24 2024, 13:30-16:30

This document contains 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) True. Since Q is an orthogonal matrix we know that $QQ^T = Q^TQ = I$. Furthermore $||A||_2 = \sqrt{\lambda_{max}(A^TA)}$. This implies:

$$\|Q^T A\|_2 = \sqrt{\lambda_{max}(A^T Q Q^T A)} = \sqrt{\lambda_{max}(A^T A)} = \|A\|_2$$

which proves the equality.

- (b) True. For the proof we use the definition of the spectral radius: $\rho(A)$ is the in absolute value largest eigenvalue of A. We know that $|\lambda| ||\mathbf{u}|| = ||A\mathbf{u}||$ for any eigenpair λ, \mathbf{u} . From the definition of a multiplicative norm ||.|| it follows that $|\lambda|||\mathbf{u}|| = ||A\mathbf{u}|| \le ||A|| ||\mathbf{u}||$. Division by $||\mathbf{u}||$ shows that $|\lambda| \le ||A||$ for any eigenvalue λ of A. So it also holds for the in absolute value largest eigenvalue of A which proves the result.
- (c) False. From $\mathbf{r} = \alpha_1 \mathbf{v}_1 + \alpha_2 \mathbf{v}_2$ it follows that $A^k \mathbf{r} = \lambda_1^k \mathbf{r}$. So all powers of A multiplied with \mathbf{r} are element of the span of \mathbf{r} . This implies that the dimension of $K^5(A, \mathbf{r})$ is equal to 1.
- (d) True. The symmetry of A_k easily follows from the symmetry of A. Since matrix A is SPD we know that $\mathbf{v}^T A \mathbf{v} > 0$ for all $\mathbf{v} \neq \mathbf{0}$. Take $\mathbf{w} \in \mathbb{R}^k$ an arbitrary vector such that $\mathbf{w} \neq \mathbf{0}$. Define $\mathbf{v} \in \mathbb{R}^n$ such that $\mathbf{v}(1) = \mathbf{w}(1), \dots, \mathbf{v}(k) = \mathbf{w}(k)$ and $\mathbf{v}(k+1) = \dots = \mathbf{v}(n) = 0$. It now follows that $\mathbf{w}^T A_k \mathbf{w} = \mathbf{v}^T A \mathbf{v} > 0$ which shows that the statement is true.
- (e) True. The condition number is defined as follows: $K_p(A) = ||A||_p ||A^{-1}||_p$. The inverse of I is equal to I. From the definition of the p-norm it follows that

$$||I||_p = \max_{||\mathbf{u}||_p=1} ||I\mathbf{u}||_p = 1$$

Combination of these expressions shows that $K_p(I) = 1$.

2. (a) The finite difference stencil is given by

$$\frac{1}{h^2}[-1\ 2 - 10h^2\ -1]$$

In order to show that the method is second order accurate, a Taylor expansion in the points x_{i-1} and x_{i+1} should be given around the point x_i where the remainder term is $O(h^4)$. It then follows that

$$-u_i'' = \frac{-u_{i-1} + 2u_i - u_{i+1}}{h^2} + O(h^2)$$

This leads to the stencil as given above.

(b) Use the goniometric formula's to show that

$$\lambda_k = -10 + \frac{2}{h^2} [1 - \cos(\pi hk)] = -10 + \frac{4}{h^2} \sin^2(\frac{\pi hk}{2})$$

(c) Since the boundary conditions are eliminated, the matrix is symmetric. Note that Gerschgorin's theorem does not imply that all eigenvalues are positive. Using the expression given in (b) with N = 10 and k = 1 shows that $\lambda_1 = -0.1314$. So eigenvalue λ_1 is negative with implies that for eigenvector \mathbf{v}_1 we have:

$$\mathbf{v}_1^T A \mathbf{v}_1 = \lambda_1 \|\mathbf{v}_1\|_2^2 < 0$$

This implies that the matrix is not SPD.

- (d) We have the following ordering:
 - 4 9 5
 7 3 8
 1 6 2
 This leads to the following non-zero pattern:

- 3. (a) Let $A\mathbf{u} = \mathbf{f}$, and A = M N where M is non singular. Derive a formula for \mathbf{u}^k and \mathbf{r}^k .
 - i.

$$\mathbf{u}^{k+1} = M^{-1}N\mathbf{u}^k + M^{-1}\mathbf{f}$$

= $M^{-1}(M-A)\mathbf{u}^k + M^{-1}\mathbf{f}$
= $\mathbf{u}^k + M^{-1}(\mathbf{f} - A\mathbf{u}^k)$
= $\mathbf{u}^k + M^{-1}\mathbf{r}^k$ (1)

ii.

$$\mathbf{r}^{k+1} = \mathbf{f} - A\mathbf{u}^{k+1}$$

= $\mathbf{f} - A(\mathbf{u}^k + M^{-1}\mathbf{r}^k)$
= $\mathbf{f} - A\mathbf{u}^k - AM^{-1}\mathbf{r}^k$
= $\mathbf{r}^k - AM^{-1}\mathbf{r}^k$
= $(I - AM^{-1})\mathbf{r}^k$ (2)

- (b) Give the iteration matrix and a sufficient condition for convergence
 - i. The iteration matrix is given by $B = I M^{-1}A$.
 - ii. There are three possible answers
 - A. $\rho(B) < 1$ B. ||B|| < 1C. $\lim_{k \to \infty} ||B^k||_2 = 0$
- (c) Assume A is upper triangular, show Jacobi converges.

Now we have M = D, where D is the matrix containing only the diagonal elements of A.

Then $B = I - D^{-1}A = I - (I + U) = U$ where U is an upper triangular matrix with zeros on the diagonal. So B is an upper triangular matrix with only zeros on the diagonal.

It then follows that $B^{n-1} = 0_{matrix}$ so the Jacobi method converges.

(d) Assume A is upper triangular, which Gauss Seidel variant is optimal for this matrix?

Solution: There are two Gauss Seidel variants: M is equal to the lower triangular part of A or M is equal to the upper triangular part of A. For this matrix the final choice is optimal. The motivation is as follows. Note that in this case M = A and therefore

$$B = I - M^{-1}A = I - A^{-1}A = 0_{matrix}.$$
(3)

Then, ||B|| < 1 and therefore this variant of GS converges. Furthermore $e^1 = Be^0 = 0_{vector}$, so the method converges after 1 iteration.

- (e) Below 3 different stopping criteria and their properties are given.
 - i. $||r^k|| \leq \epsilon$, this criterion is not scaling invariant.
 - ii. $\frac{\|r^k\|}{\|r^0\|} \leq \epsilon$, depends on the quality of the initial guess.
 - iii. $\frac{\|r^k\|}{\|f\|} \leq \epsilon$, this is a good stopping criterion.
- 4. (a) We take $\mathbf{u}^1 = \alpha \mathbf{r}^0$ where α is a constant which has to be chosen such that $\|\mathbf{u} \mathbf{u}^1\|_2$ is minimal. This leads to

$$\|\mathbf{u} - \mathbf{u}^1\|_2^2 = (\mathbf{u} - \alpha \mathbf{r}^0)^T (\mathbf{u} - \alpha \mathbf{r}^0) = \mathbf{u}^T \mathbf{u} - 2\alpha (\mathbf{r}^0)^T \mathbf{u} + \alpha^2 (\mathbf{r}^0)^T \mathbf{r}^0.$$

The norm given above is minimized if $\alpha = \frac{(\mathbf{r}^0)^T \mathbf{u}}{(\mathbf{r}^0)^T \mathbf{r}^0}$.

(b) Note that **u** is needed in the definition of α as given in part (a). If one uses the A-norm, the following expression should be minimal: $\|\mathbf{u} - \mathbf{u}^1\|_A$. Using the same steps as in part (a) it appears that now α is given by:

$$\alpha = \frac{(\mathbf{r}^0)^T A \mathbf{u}}{(\mathbf{r}^0)^T A \mathbf{r}^0} = \frac{(\mathbf{r}^0)^T \mathbf{f}}{(\mathbf{r}^0)^T A \mathbf{r}^0},$$

which is easy to compute.

(c) The optimality property of CG implies that the approximation \mathbf{u}^k coming from CG satisfies:

$$\|\mathbf{u} - \mathbf{u}^k\|_A = \min_{\mathbf{y} \in K^k(A; \mathbf{r}^0)} \|\mathbf{u} - \mathbf{y}\|_A$$

If the method terminates before we reach k = n we know that we have a 'lucky' breakdown so $\mathbf{u}^k = \mathbf{u}$. If not we know that the dimension of K^n is equal to n, thus $K^n = \mathbb{R}^n$ and thus $\mathbf{u}^n = \mathbf{u}$.

(d) 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 $K_2(A_1) = 10$ and $K_2(A_2) = 200$, it is clear that we expect that the convergence for A_1 is much faster than for A_2 .

- (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 low cost.
- 5. (a) If we do the multiplication:

$$(I - \alpha^{(k)} \mathbf{e}_k^T) (I + \alpha^{(k)} \mathbf{e}_k^T)$$

we obtain the following:

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

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

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

which proves the claim that $M_k^{-1} = I + \alpha^{(k)} \mathbf{e}_k^T$.

(b) The perturbed solution $\mathbf{u} + \Delta \mathbf{u}$ solves the system

$$A(\mathbf{u} + \Delta \mathbf{u}) = \mathbf{f} + \Delta \mathbf{f} \,. \tag{4}$$

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

$$A\Delta \mathbf{u} = \Delta \mathbf{f} \,, \tag{5}$$

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

$$\frac{1}{\|\mathbf{u}\|} \le \|A\| \ \frac{1}{\|\mathbf{f}\|} \tag{6}$$

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

$$\frac{\|\Delta \mathbf{u}\|}{\|\mathbf{u}\|} \le \|A^{-1}\| \|A\| \frac{\|\Delta \mathbf{f}\|}{\|\mathbf{f}\|} = \kappa(A) \frac{\|\Delta \mathbf{f}\|}{\|\mathbf{f}\|} \le \delta \kappa(A),$$
(7)

where $\kappa(A)$ denotes the condition number of A measured in the norm $\|\cdot\|$.

(c) The LU decomposition determines an upper triangular matrix U and a lower triangular matrix L, with $l_{ii} = 1$, where A = LU. The procedure to obtain this decomposition is using Gauss transformations, such that column k is transformed in a such a way that all element k + 1, ..., n of this column become equal to zero. Since the non-zero pattern of L and U is the same as that of A, the number of computations per row are: one division to compute the multiplier followed by a multiplication and addition/subtraction to compute the diagonal element of U. This leads to 3n flops for the decomposition.

In order to find solution \mathbf{u} from $A\mathbf{u} = \mathbf{f}$, we substitute the decomposition into $A\mathbf{u} = \mathbf{f}$, so $LU\mathbf{u} = \mathbf{f}$. If we define $\mathbf{y} = U\mathbf{u}$ we can first solve $L\mathbf{y} = \mathbf{f}$ and then $U\mathbf{u} = \mathbf{y}$. Since these systems are both triangular this is easy to solve. The work for the first system is one multiplication and addition/subtraction to compute one component of \mathbf{y} , so in total 2n flops for the first system. For the second system one needs a multiplication and addition/subtraction followed by a division to compute one component of \mathbf{u} , so in total 3n flops for the second system. In total 8n flops are needed to solve a system with a tri-diagonal matrix A.

(d) After 1 step of the Gaussian elimination process we obtain the following matrix:

$$\begin{pmatrix} 4 & -1 & 0 & 0 & -1 \\ 0 & 3\frac{3}{4} & -1 & 0 & -\frac{1}{4} \\ 0 & -1 & 4 & -1 & 0 \\ 0 & 0 & -1 & 4 & -1 \\ 0 & -\frac{1}{4} & 0 & -1 & 3\frac{3}{4} \end{pmatrix}$$

Note that the fill in less than $\frac{1}{4}$ in absolute value.