

Chapter 3

Classical Iterative Schemes

3.1 General Theory

Remark 3.1. Basic idea, transform to a fixed-point equation. The construction of a classical iterative scheme for solving (1.1) starts with the decomposition

$$A = M - N, \quad M, N \in \mathbb{R}^{n \times n}, \quad M \text{ is non-singular,}$$

of the system matrix A . Using this decomposition, (1.1) can be transformed into the fixed-point equation

$$M\underline{x} = \underline{b} + N\underline{x} \iff \underline{x} = M^{-1}(\underline{b} + N\underline{x}). \quad (3.1)$$

Given an initial iterate $\underline{x}^{(0)} \in \mathbb{R}^n$, one can try to solve (3.1) with the fixed-point iteration

$$\underline{x}^{(k+1)} = M^{-1}(\underline{b} + N\underline{x}^{(k)}), \quad k = 0, 1, 2, \dots \quad (3.2)$$

Banach's¹ fixed-point theorem gives information on the convergence of this iteration. □

Theorem 3.2. Banach's fixed-point theorem. *Let (\mathcal{X}, d) be a complete metric space and let $f : \mathcal{X} \rightarrow \mathcal{X}$ be a contraction (f is Lipschitz² continuous with the Lipschitz constant $L < 1$). Then, the equation $x = f(x)$ possesses a unique solution $\hat{x} \in \mathcal{X}$ (a fixed point). The iterative scheme*

$$x^{(k+1)} = f(x^{(k)}), \quad k = 0, 1, 2, \dots$$

converges to \hat{x} for any initial iterate $x^{(0)} \in \mathcal{X}$.

Proof. Basic course on calculus. ■

Theorem 3.3. Condition on the iteration matrix of (3.2) for convergence. *The iterative scheme (3.2) converges to the solution \underline{x} of (1.1) for any initial iterate $\underline{x}^{(0)}$ if and only if the spectral radius of the iteration matrix $G = M^{-1}N$ is smaller than one: $\rho(G) < 1$.*

Proof. *i) Lipschitz continuity.* The iteration (3.2) is a fixed-point iteration with

$$f : \mathbb{R}^n \rightarrow \mathbb{R}^n, \quad \underline{x} \mapsto M^{-1}N\underline{x} + M^{-1}\underline{b}.$$

Let \mathbb{R}^n be equipped with the vector norm $\|\cdot\|_*$ defined in (2.2) with sufficiently small δ . The operator $G = M^{-1}N$ is linear and bounded since $\|G\|_*$ is finite. Hence, G is continuous and even Lipschitz continuous. Since f is continuously differentiable, the Lipschitz constant is given by

¹ Stefan Banach (1892 – 1945)

² Rudolf Lipschitz (1832 – 1903)

$$L_* = \sup_{\underline{x} \in \mathbb{R}^n} \|J(f(\underline{x}))\|_* = \sup_{\underline{x} \in \mathbb{R}^n} \|G\|_* = \|G\|_*,$$

where $J(f(\underline{x}))$ is the (constant) Jacobian of $f(\underline{x})$.

ii) *Convergence for $\rho(G) < 1$.* Let $\rho(G) < 1$. Then, it is possible to find a matrix norm $\|\cdot\|_*$ such that, according to Lemma 2.8, $\|G\|_* \leq \rho(G) + \varepsilon < 1$ with $\varepsilon > 0$. Hence, $L_* < 1$ and $f(\underline{x})$ is a contraction. Now, the statement follows with Theorem 3.2.

iii) $\rho(G) \geq 1$. Let $\rho(G) \geq 1$. An initial guess will be constructed for which the fixed-point iteration does not converge. Without loss of generality, consider the case $\underline{b} = \underline{0}$ such that the solution of (1.1) is $\underline{x} = \underline{0}$.

Since $\rho(G) \geq 1$, there is an eigenvalue $\lambda \in \mathbb{C}$ of G with $|\lambda| \geq 1$. This eigenvalue can be written in the form

$$\lambda = |\lambda| (\cos(\varphi) + i \sin(\varphi)), \quad (3.3)$$

where φ is the argument of λ . Let $\underline{z} \in \mathbb{C}^n$, $\underline{z} \neq \underline{0}$, be a corresponding eigenvector:

$$G\underline{z} = \lambda\underline{z}. \quad (3.4)$$

From the conjugate of this equation $\overline{G\underline{z}} = \overline{\lambda\underline{z}}$, it follows that $G\overline{\underline{z}} = \overline{\lambda}\overline{\underline{z}}$ since G is a real matrix.

Choose the initial iterate $\underline{x}^{(0)} = \underline{z} + \overline{\underline{z}} \in \mathbb{R}^n$.

One has to exclude that $\underline{x}^{(0)} = \underline{0}$. If $\underline{x}^{(0)} = \underline{0}$, then $\underline{z} = i\underline{v}$ with $\underline{v} \in \mathbb{R}^n$. One obtains from the eigenvalue equation that $iG\underline{v} = i\lambda\underline{v}$ which is equivalent to $G\underline{v} = \lambda\underline{v}$. On the left-hand side of this equation, there is a real vector. Since \underline{v} is a real vector, it follows that λ must be real, too. But in this case, the corresponding eigenvector \underline{z} is also real and it cannot be of the form $\underline{z} = i\underline{v}$. Hence, an eigenvector of form $\underline{z} = i\underline{v}$ cannot exist and $\underline{x}^{(0)} \neq \underline{0}$.

Using the definition of the initial iterate, the eigenvalue problem (3.4), and some basic properties of eigenvalues and complex numbers, it follows that

$$\underline{x}^{(k)} = \underbrace{G \left(G \left(\dots G \underline{x}^{(0)} \right) \right)}_{k \text{ times}} = G^k \underline{x}^{(0)} = G^k \underline{z} + G^k \overline{\underline{z}} = \lambda^k \underline{z} + \overline{\lambda}^k \overline{\underline{z}} = 2\operatorname{Re} \left(\lambda^k \underline{z} \right), \quad k = 0, 1, \dots$$

The iteration converges to the solution $\underline{x} = \underline{0}$ if

$$\begin{aligned} \underline{0} &= \lim_{k \rightarrow \infty} 2\operatorname{Re} \left(\lambda^k \underline{z} \right) = \lim_{k \rightarrow \infty} 2|\lambda|^k \operatorname{Re} \left((\cos(k\varphi) + i \sin(k\varphi)) \underline{z} \right) \\ &= \lim_{k \rightarrow \infty} 2|\lambda|^k (\cos(k\varphi) \operatorname{Re}(\underline{z}) - \sin(k\varphi) \operatorname{Im}(\underline{z})), \end{aligned}$$

where (3.3) and basic properties of the real part of complex numbers were used. The factor $|\lambda|^k$ is always larger or equal to 1, since $|\lambda| \geq 1$. Hence, the second factor has to converge to zero if the iteration should converge to $\underline{x} = \underline{0}$. Note that the second factor is a vector. It converges to zero if and only if each of its components converges to zero. There is at least one component z_l with $z_l \neq 0$, since \underline{z} is an eigenvector. Let ζ be the argument of z_l , i.e., the angle in the complex plane. Using the definition of the real and imaginary part of z_l and the trigonometric identity for $\cos(\alpha + \beta)$ yields

$$\cos(k\varphi) \operatorname{Re}(z_l) - \sin(k\varphi) \operatorname{Im}(z_l) = |z_l| (\cos(k\varphi) \cos(\zeta) - \sin(k\varphi) \sin(\zeta)) = |z_l| \cos(k\varphi + \zeta).$$

The only possibility to obtain convergence to zero for $k \rightarrow \infty$ for the periodic cosine function and for fixed increment φ is the case $\zeta = \pm\pi/2$ and $\varphi \in \{0, \pi\}$, because then $k\varphi + \zeta$ is $\pi/2$ plus an integer multiple of π , so that the term with the cosine vanishes. Since $\zeta = \pm\pi/2$, it follows that $\operatorname{Re}(z_l) = 0$ and since $\varphi \in \{0, \pi\}$, it is $\lambda \in \mathbb{R}$. However, if $\lambda \in \mathbb{R}$, then the eigenvector \underline{z} is real, too, which is a contradiction to $\operatorname{Re}(z_l) = 0$ and $z_l \neq 0$.

In summary, the iteration (3.2) does not converge for the initial iterate $\underline{x}^{(0)} = \underline{z} + \overline{\underline{z}}$. That means, if the iteration (3.2) converges for all initial iterates, then $\rho(G) \geq 1$ cannot hold. \blacksquare

Remark 3.4. Complex-valued systems. The last part of the proof simplifies much if one considers complex-valued systems of linear equations. Then, one can take the initial iterate $\underline{x}^{(0)} = \underline{z} \neq \underline{0}$, finds that $\underline{x}^{(k)} = \lambda^k \underline{z}$, and concludes that $\|\underline{x}^{(k)}\|_2 = |\lambda|^k \|\underline{z}\|_2 \not\rightarrow 0$, since $\|\underline{z}\|_2 \neq 0$ and $|\lambda|^k \geq 1$.

However, if a real-valued system is given, computations in practice are performed usually only with real-valued vectors. Thus, the goal of the proof was to show that there is a real-valued initial iterate for which the iteration does not converge. Showing that there is a complex-valued initial iteration for which the iteration does not converge does not exclude that it converges for all real-valued initial iterates. In particular, the real-valued situation is not a special case of the complex-valued case. \square

3.2 Examples for Classical Iterative Schemes

Remark 3.5. Decomposition of the system matrix. One uses for the definition of classical iterative schemes a decomposition of the matrix A in the form

$$A = D + L + U,$$

where D is the diagonal of A , L is its strict lower part and U its strict upper part.

It turns out that the presented classical iterative schemes require that the diagonal of A does not contain a zero entry. \square

Example 3.6. Jacobi method. The Jacobi method is derived by setting

$$M = D, \quad N = -(L + U).$$

A straightforward calculation reveals that the fixed-point equation (3.1) has the form

$$\underline{x} = D^{-1}(\underline{b} - (L + U)\underline{x}) = \underline{x} + D^{-1}(\underline{b} - A\underline{x}).$$

This fixed-point equation gives the following iterative scheme, called Jacobi³ method

$$\underline{x}^{(k+1)} = \underline{x}^{(k)} + D^{-1}(\underline{b} - A\underline{x}^{(k)}), \quad k = 0, 1, 2, \dots \quad (3.5)$$

The iteration matrix is $G_{\text{Jac}} = -D^{-1}(L + U)$. \square

Example 3.7. Damped Jacobi method. Let $\omega \in \mathbb{R}$, $\omega > 0$. The matrices that define the fixed-point equation for the damped Jacobi method are given by

$$M = \omega^{-1}D, \quad N = \omega^{-1}D - A.$$

The damped Jacobi method has the form

$$\underline{x}^{(k+1)} = \underline{x}^{(k)} + \omega D^{-1}(\underline{b} - A\underline{x}^{(k)}), \quad k = 0, 1, 2, \dots \quad (3.6)$$

and the iteration matrix is $G_{\text{dJac}} = I - \omega D^{-1}A$.

Damping is a very important concept for all iterative methods (for linear or nonlinear problems). Comparing (3.5) and (3.6), one can see that for $\omega < 1$ the damped Jacobi method does not accept the whole update proposed by the Jacobi method, but only a part of this update. In practice, damping might be very useful to achieve convergence. But there are situations where it does not help, compare also Example 3.13. \square

Example 3.8. Gauss–Seidel method. In the Gauss⁴–Seidel⁵ method, the invertible matrix M is a triangular matrix

$$M = D + L, \quad N = -U.$$

It follows that

$$\begin{aligned} \underline{x}^{(k+1)} &= (D + L)^{-1}(\underline{b} - U\underline{x}^{(k)}) \\ &= (D + L)^{-1}(\underline{b} - A\underline{x}^{(k)} + (A - U)\underline{x}^{(k)}) \end{aligned}$$

³ Carl Gustav Jacob Jacobi (1804 – 1851)

⁴ Johann Carl Friedrich Gauss (1777 – 1855)

⁵ Philipp Ludwig von Seidel (1821 – 1896)

$$= \underline{x}^{(k)} + (D + L)^{-1} \left(\underline{b} - A\underline{x}^{(k)} \right), \quad k = 0, 1, 2, \dots$$

The iteration matrix has the form $G_{\text{GS}} = M^{-1}N = -(D + L)^{-1}U$. Multiplying the first equation of the Gauss–Seidel method from left by $(D + L)$ and then by D^{-1} and rearranging terms gives the more familiar form of this iteration

$$\begin{aligned} \underline{x}^{(k+1)} &= D^{-1} \left(\underline{b} - L\underline{x}^{(k+1)} - U\underline{x}^{(k)} \right) \\ &= \underline{x}^{(k)} + D^{-1} \left(\underline{b} - L\underline{x}^{(k+1)} - (D + U)\underline{x}^{(k)} \right), \quad k = 0, 1, 2, \dots \end{aligned}$$

Writing this iteration for the components of the vector shows that the right-hand side can be evaluated even if the new iterate appears there, since only already computed components of the new iterate are needed for this evaluation

$$x_i^{(k+1)} = x_i^{(k)} + \frac{1}{a_{ii}} \left(b_i - \sum_{j=1}^{i-1} a_{ij}x_j^{(k+1)} - \sum_{j=i}^n a_{ij}x_j^{(k)} \right), \quad k = 0, 1, 2, \dots$$

□

Example 3.9. SOR method. The matrices that define the (forward) successive over relaxation (SOR) method are given by

$$M = \omega^{-1}D + L, \quad N = \omega^{-1}D - (D + U),$$

where $\omega \in \mathbb{R}, \omega > 0$. This method can be written in the form

$$\underline{x}^{(k+1)} = \underline{x}^{(k)} + \omega D^{-1} \left(\underline{b} - L\underline{x}^{(k+1)} - (D + U)\underline{x}^{(k)} \right), \quad k = 0, 1, 2, \dots$$

or

$$\underline{x}^{(k+1)} = \underline{x}^{(k)} + \left(\frac{D}{\omega} + L \right)^{-1} \left(\underline{b} - A\underline{x}^{(k)} \right), \quad k = 0, 1, 2, \dots \quad (3.7)$$

For $\omega = 1$, the Gauss–Seidel method is recovered. One obtains for the iteration matrix

$$\begin{aligned} G_{\text{SOR}}(\omega) &= \left(\omega^{-1}D + L \right)^{-1} \left(\omega^{-1}D - (D + U) \right) \\ &= \omega (D + \omega L)^{-1} \left(\omega^{-1}D - (D + U) \right) \\ &= (D + \omega L)^{-1} ((1 - \omega)D - \omega U). \end{aligned} \quad (3.8)$$

□

Example 3.10. SSOR method. In the SOR method, one can change the roles of L and U to obtain the backward SOR method

$$\underline{x}^{(k+1)} = \underline{x}^{(k)} + \omega D^{-1} \left(\underline{b} - U\underline{x}^{(k+1)} - (D + L)\underline{x}^{(k)} \right), \quad k = 0, 1, 2, \dots$$

This method updates the unknowns in reverse order. The forward and backward SOR behave in general differently. There are cases where one of them works much more efficient than the other one. However, in general one does not know a priori which is the better variant. The SSOR (symmetric SOR) method combines both methods. One step of SSOR consists of two substeps, one forward SOR and one backward SOR step:

$$\begin{aligned} \underline{x}^{(k+1/2)} &= \underline{x}^{(k)} + \omega D^{-1} \left(\underline{b} - L\underline{x}^{(k+1/2)} - (D + U)\underline{x}^{(k)} \right) \\ \underline{x}^{(k+1)} &= \underline{x}^{(k+1/2)} + \omega D^{-1} \left(\underline{b} - U\underline{x}^{(k+1)} - (D + L)\underline{x}^{(k+1/2)} \right), \quad k = 0, 1, 2, \dots \end{aligned}$$

□

3.3 Some Convergence Results

Theorem 3.11. Convergence for strongly diagonally dominant matrices. *Let $A \in \mathbb{R}^{n \times n}$ be strongly diagonally dominant. Then, the Jacobi method and the Gauss–Seidel method converge for every initial iterate $\underline{x}^{(0)} \in \mathbb{R}^n$.*

Proof. Following Theorem 3.3, one has to show that the spectral radius of the iteration matrices is smaller than 1. *Jacobi method.* Let $\underline{z} \in \mathbb{R}^n, \underline{z} \neq \underline{0}$. Then, the triangle inequality gives

$$\begin{aligned} |(G_{\text{Jac}}\underline{z})_i| &= \left| \left(-D^{-1}(L+U)\underline{z} \right)_i \right| = \left| \frac{1}{a_{ii}} \sum_{j=1, j \neq i}^n a_{ij} z_j \right| \leq \frac{1}{|a_{ii}|} \sum_{j=1, j \neq i}^n |a_{ij}| |z_j| \\ &\leq \frac{1}{|a_{ii}|} \underbrace{\sum_{j=1, j \neq i}^n |a_{ij}|}_{< |a_{ii}|} \|\underline{z}\|_\infty < \|\underline{z}\|_\infty, \quad i = 1, \dots, n, \end{aligned}$$

which is equivalent to $\|G_{\text{Jac}}\underline{z}\|_\infty < \|\underline{z}\|_\infty$. With Remark 2.7 and the definition of $\|\cdot\|_\infty$, Remark 2.3, it follows that

$$\rho(G_{\text{Jac}}) \leq \|G_{\text{Jac}}\|_\infty = \max_{\underline{z} \in \mathbb{R}^n, \underline{z} \neq \underline{0}} \frac{\|G_{\text{Jac}}\underline{z}\|_\infty}{\|\underline{z}\|_\infty} < 1. \quad (3.9)$$

Gauss–Seidel method. A direct calculation shows (exercise)

$$G_{\text{GS}} = -D^{-1}(LG_{\text{GS}} + U). \quad (3.10)$$

This relation and the triangle inequality gives for the first component and $\underline{z} \in \mathbb{R}^n, \underline{z} \neq \underline{0}$,

$$|(G_{\text{GS}}\underline{z})_1| \leq \frac{1}{|a_{11}|} \sum_{j=2}^n |a_{1j}| |z_j| \leq \frac{1}{|a_{11}|} \underbrace{\sum_{j=2}^n |a_{1j}|}_{< |a_{11}|} \|\underline{z}\|_\infty < \|\underline{z}\|_\infty,$$

where the term with the factor LG_{GS} vanishes since the first row of L consists only of zeros. Using now (3.10) and the induction $|(G_{\text{GS}}\underline{z})_j| < \|\underline{z}\|_\infty, j < i$ yields

$$\begin{aligned} |(G_{\text{GS}}\underline{z})_i| &\leq \frac{1}{|a_{ii}|} \left(\sum_{j=1}^{i-1} |a_{ij}| |(G_{\text{GS}}\underline{z})_j| + \sum_{j=i+1}^n |a_{ij}| |z_j| \right) \\ &\leq \frac{1}{|a_{ii}|} \underbrace{\left(\sum_{j=1, j \neq i}^n |a_{ij}| \right)}_{< |a_{ii}|} \|\underline{z}\|_\infty < \|\underline{z}\|_\infty, \quad i = 2, \dots, n. \end{aligned}$$

The remainder of the proof is like for the Jacobi method, compare (3.9). ■

Lemma 3.12. Eigenvalues of the iteration matrix of the damped Jacobi method. *Let $\omega > 0$, then $\lambda \in \mathbb{C}$ is an eigenvalue of G_{Jac} if and only if $\mu = 1 - \omega + \omega\lambda$ is an eigenvalue of G_{dJac} .*

Proof. It is with $A = D + L + U$

$$G_{\text{dJac}} = I - \omega D^{-1}A = I - \omega D^{-1}D - \omega \underbrace{D^{-1}(L+U)}_{-G_{\text{Jac}}} = (1 - \omega)I + \omega G_{\text{Jac}}.$$

The statement of the lemma follows now from well known properties of eigenvalues. ■

Example 3.13. Convergence of the damped Jacobi method where the Jacobi method fails. If ω is chosen appropriately, there is the possibility that the damped Jacobi method converges for every initial guess whereas the Jacobi method does not.

Assume that G_{Jac} has only real eigenvalues. Denote by λ_{\min} the smallest one and by λ_{\max} the largest one. If

$$\lambda_{\min} < -1 < \lambda_{\max} < 1,$$

then there are initial iterates for which the Jacobi method does not converge, compare Theorem 3.3. From Lemma 3.12, one has

$$\begin{aligned}\mu_{\min} &= (1 - \omega) + \omega\lambda_{\min} = 1 - \omega(1 - \lambda_{\min}), \\ \mu_{\max} &= (1 - \omega) + \omega\lambda_{\max} = 1 - \omega(1 - \lambda_{\max}).\end{aligned}\tag{3.11}$$

It follows that

$$-1 < \mu_{\min} < 1 \quad \text{if } \omega < \frac{2}{1 - \lambda_{\min}} < 1, \quad -1 < \mu_{\max} < 1 \quad \text{if } 0 < \omega \leq 1.$$

The choice $\omega \in (0, 2/(1 - \lambda_{\min}))$ ensures the convergence of the damped Jacobi method for each initial iterate.

Consider the case $\lambda_{\max} > 1$. Then, one gets from (3.11) that $\mu_{\max} > 1$. In this case, there are initial iterates for which the damped Jacobi method does not converge as well. \square

Lemma 3.14. Parameter in the case that the SOR method converges, Lemma of Kahan⁶. *If the SOR method converges for every initial iterate $\underline{x}^{(0)} \in \mathbb{R}^n$, then $\omega \in (0, 2)$.*

Proof. Let $\lambda_1, \dots, \lambda_n \in \mathbb{C}$ be the eigenvalues of $G_{\text{SOR}}(\omega)$. It is with (3.8) and properties of the determinant

$$\begin{aligned}\prod_{i=1}^n \lambda_i &= \det(G_{\text{SOR}}(\omega)) = \det\left((D + \omega L)^{-1}((1 - \omega)D - \omega U)\right) \\ &= \det\left(\underbrace{(D + \omega L)^{-1}}_{\text{lower triangular matrix}}\right) \det\left(\underbrace{(1 - \omega)D - \omega U}_{\text{upper triangular matrix}}\right) \\ &= \det(D^{-1})(1 - \omega)^n \det(D) = (1 - \omega)^n.\end{aligned}$$

Hence, it follows that

$$\prod_{i=1}^n |\lambda_i| = |1 - \omega|^n.$$

There is at least one eigenvalue λ_i with $|\lambda_i| \geq |1 - \omega|$ and it follows that $\rho(G_{\text{SOR}}(\omega)) \geq |1 - \omega|$. Now, the application of Theorem 3.3 shows that SOR cannot converge for all initial iterates if $\omega \notin (0, 2)$, because if $\omega \notin (0, 2)$ then $\rho(G_{\text{SOR}}(\omega)) \geq |1 - \omega| \geq 1$. \blacksquare

Theorem 3.15. Convergence of SOR for s.p.d. matrices. *Let $A \in \mathbb{R}^{n \times n}$ be s.p.d. Then the SOR method converges for all initial iterates $\underline{x}^{(0)} \in \mathbb{R}^n$ if $\omega \in (0, 2)$.*

Proof. For interested students, the special case $\omega = 1$ will be an exercise problem. Let $\lambda \in \mathbb{C}$ be an eigenvalue of $G_{\text{SOR}}(\omega)$, see (3.8), and let $\underline{z} \in \mathbb{C}^n$ be a corresponding eigenvector, i.e.,

$$(D + \omega L)^{-1}((1 - \omega)D - \omega U)\underline{z} = \lambda\underline{z}.\tag{3.12}$$

Following Theorem 3.3, one has to find a condition such that $|\lambda| < 1$ is satisfied. The following identities can be easily verified using the decomposition of A :

⁶ William M. Kahan, born 1933

$$\begin{aligned} D + \omega L &= \left(1 - \frac{\omega}{2}\right) D + \frac{\omega}{2} A + \frac{\omega}{2} (L - U), \\ (1 - \omega) D - \omega U &= \left(1 - \frac{\omega}{2}\right) D - \frac{\omega}{2} A + \frac{\omega}{2} (L - U). \end{aligned}$$

Inserting these identities in the eigenvalue equation (3.12) and multiplying this equation from the left-hand side with $(D + \omega L)$ and with the adjoint vector \underline{z}^* , one obtains

$$\lambda = \frac{\left(1 - \frac{\omega}{2}\right) \underline{z}^* D \underline{z} - \frac{\omega}{2} \underline{z}^* A \underline{z} + \frac{\omega}{2} \underline{z}^* (L - U) \underline{z}}{\left(1 - \frac{\omega}{2}\right) \underline{z}^* D \underline{z} + \frac{\omega}{2} \underline{z}^* A \underline{z} + \frac{\omega}{2} \underline{z}^* (L - U) \underline{z}}.$$

Now, the terms in this expression will be considered individually. The matrix $L - U$ is skew-symmetric since A is symmetric. It follows for all $\underline{x} \in \mathbb{R}^n$ that

$$\underbrace{\underline{x}^T (L - U) \underline{x}}_{\in \mathbb{R}} = \left(\underline{x}^T (L - U) \underline{x}\right)^T = \underline{x}^T (L - U)^T \underline{x} = -\underline{x}^T (L - U) \underline{x} \in \mathbb{R},$$

consequently $\underline{x}^T (L - U) \underline{x} = 0$ for all $\underline{x} \in \mathbb{R}^n$, and

$$\operatorname{Re}(\underline{z}^* (L - U) \underline{z}) = \underbrace{\operatorname{Re}(\underline{z}^*)}_{\in \mathbb{R}^n} (L - U) \underbrace{\operatorname{Re}(\underline{z})}_{\in \mathbb{R}^n} + \underbrace{\operatorname{Im}(\underline{z}^*)}_{\in \mathbb{R}^n} (L - U) \underbrace{\operatorname{Im}(\underline{z})}_{\in \mathbb{R}^n} = 0.$$

Hence, $\underline{z}^* (L - U) \underline{z} = ia$ with $a \in \mathbb{R}$.

Since A is positive definite, its diagonal D is positive definite, too. The products $\underline{z}^* D \underline{z}$ and $\underline{z}^* A \underline{z}$ are positive real numbers since for $\underline{z} = \underline{u} + i\underline{v}$, $\underline{z}^* = \underline{u}^T - i\underline{v}^T$, $\underline{u}, \underline{v} \in \mathbb{R}^n$, $\underline{z} \neq \underline{0}$ ($\underline{u} \neq \underline{0}$ or $\underline{v} \neq \underline{0}$), because it is an eigenvector, one obtains with the symmetry of A and the positive definiteness of this matrix

$$\begin{aligned} \underline{z}^* A \underline{z} &= \underline{u}^T A \underline{u} - i\underline{v}^T A \underline{u} + i\underline{u}^T A \underline{v} - i^2 \underline{v}^T A \underline{v} \\ &= \underline{u}^T A \underline{u} - i\underline{v}^T A \underline{u} + i\underline{u}^T A \underline{v} + \underline{v}^T A \underline{v} > 0. \end{aligned}$$

It follows that λ has the form

$$\lambda = \frac{b + ia}{c + ia} \quad a, b, c \in \mathbb{R}, c > 0,$$

so that the denominator is not zero, with

$$b = \left(1 - \frac{\omega}{2}\right) \underline{z}^* D \underline{z} - \frac{\omega}{2} \underline{z}^* A \underline{z}, \quad c = \left(1 - \frac{\omega}{2}\right) \underline{z}^* D \underline{z} + \frac{\omega}{2} \underline{z}^* A \underline{z}.$$

Consequently, it is

$$|\lambda|^2 = \frac{b^2 + a^2}{c^2 + a^2} = \frac{\left[\left(1 - \frac{\omega}{2}\right) \underline{z}^* D \underline{z} - \frac{\omega}{2} \underline{z}^* A \underline{z}\right]^2 + a^2}{\left[\left(1 - \frac{\omega}{2}\right) \underline{z}^* D \underline{z} + \frac{\omega}{2} \underline{z}^* A \underline{z}\right]^2 + a^2}.$$

Thus $|\lambda| < 1$ holds only if the numerator is smaller than the denominator. The only difference of the numerator and the denominator is the mixed term in the square. Applying the square and neglecting the same positive terms on both side, this condition is equivalent to

$$\begin{aligned} - \underbrace{\omega}_{>0} \left(1 - \frac{\omega}{2}\right) \underbrace{\underline{z}^* D \underline{z}}_{>0} \underbrace{\underline{z}^* A \underline{z}}_{>0} &< \omega \left(1 - \frac{\omega}{2}\right) \underline{z}^* D \underline{z} \underline{z}^* A \underline{z} \iff \\ - \left(1 - \frac{\omega}{2}\right) &< \left(1 - \frac{\omega}{2}\right) \iff \\ \omega &< 2. \end{aligned}$$

Hence, the SOR method converges for all initial iterates if $\omega \in (0, 2)$. ■

Remark 3.16. Difficulty of choosing ω in practice. For choosing ω such that the SOR method converges as fast as possible, one needs information about the eigenvalues of A . However, the computation of these information is very costly, see Numerical Mathematics I or the literature. □

Remark 3.17. Number of iterations in practice. If classical iterative schemes are used for the solution of linear systems of equations that arise in discretizing partial differential equations, one finds that the number of

iterations to fulfill a certain stopping criterion increases rapidly if the mesh is refined. One can show that the number of iteration depends on the condition number of the matrices and it scales linearly with the condition number. As example, the standard finite difference discretization of the Poisson problem (Laplace equation) on an equidistant grid of size h leads to matrices with a condition number of $\mathcal{O}(h^{-2})$, see homework problems. It follows that the number of iterations for the solution of the linear system increases approximately by the factor 4 with each refinement $h \rightarrow h/2$. For this reason, the classical iterative schemes are not useful as solver for such systems. They are important as preconditioners or as smoothers in multigrid methods, see Chapter 8. □