# The Optimal Relaxation Parameter for the SOR Method Applied to a Classical Model Problem 

Shiming Yang* and Matthias K. Gobbert*


#### Abstract

The successive overrelaxation (SOR) method is an example of a classical iterative method for the approximate solution of a system of linear equations. Its iteration matrix depends on a relaxation parameter. There is no explicit formula for the optimal relaxation parameter in terms of properties of the system matrix of a general system matrix. However, for the classical model problem of a finite difference approximation to the Poisson equation, a formula for the optimal relaxation parameter can be derived. Beyond this model problem, this result is also useful as guidance for the choice of the parameter in other problems. This paper presents the detailed derivation of the formula for the optimal relaxation parameter for the model problem and extends the well-known one- and two-dimensional results to higher dimensions.


Key words. Iterative methods, SOR, Optimal relaxation parameter, Poisson equation, Finite differences.

AMS subject classifications (2000). 15A18, 15A90, 65F10, 65N22, 65N25.

## 1 Introduction

A classical model problem for the performance comparison of linear solvers is the system of linear equations resulting from the finite difference discretization of the Poisson equation with Dirichlet boundary conditions

$$
\begin{align*}
-\Delta u & =f \quad \text { in } \Omega,  \tag{1.1}\\
u & =g \quad \text { on } \partial \Omega, \tag{1.2}
\end{align*}
$$

see, e.g., [1, Section 6.3], [3], [4, Subsection 9.1.1], [6, Chapter 10], [11, Section 7.1]. Here, $\Omega \subset \mathbb{R}^{d}$ denotes the domain of the partial differential equation in $d=1,2,3, \ldots$ dimensions and $\partial \Omega$ denotes the boundary of $\Omega$. The Laplace operator $\Delta u$ in $d$ dimensions is defined as $\triangle u=\sum_{i=1}^{d} \frac{\partial^{2} u}{\partial x_{i}^{2}}$. This discretization results in a system of linear equations with a symmetric positive definite system matrix, for which the convergence of several classical iterative methods (Jacobi, Gauss-Seidel, SOR, SSOR) and modern iterative methods (conjugate gradient) can be guaranteed. Therefore, it is a good model problem to compare the efficiency of these and other linear solvers. The two-dimensional version of the problem, $\Omega \subset \mathbb{R}^{2}$, is the most popular example, because direct methods (Gaussian elimination) and iterative methods are competitive with each other on today's computers; by contrast, for the onedimensional version, direct solvers are optimal, and for the three-dimensional version, only iterative methods allow the solution of problems with reasonably fine spatial resolutions.

The classical iterative methods include the successive overrelaxation (SOR) method, whose formulation depends on a relaxation parameter $\omega$. The speed of convergence of classical iterative methods

[^0]is determined by the spectral radius $\rho(G)$, defined as the absolute value of the largest eigenvalue in magnitude, of the iteration matrix $G$. The value of $\omega$ for which the spectral radius is minimized and the convergence of the SOR method is the fastest is $\omega_{\text {opt }}$. For a general linear system $A x=b$, no explicit formula exists to compute its value in terms of properties of the system matrix $A$. But for the model problem (1.1)-(1.2) discretized by a finite difference method on a uniform mesh with $N+2$ points in each coordinate direction, the formula
\[

$$
\begin{equation*}
\omega_{o p t}=\frac{2}{1+\sin \pi h} \tag{1.3}
\end{equation*}
$$

\]

gives the optimal value for $\omega$ in terms of the mesh spacing $h=1 /(N+1)$.
The exact formula in (1.3) is available in several sources, e.g., [11, page 540], or sometimes an approximation such as $\omega_{\text {opt }}=2(1-\pi h)+\mathcal{O}\left(h^{2}\right)$ obtained using a Taylor expansion for the sine function [4, page 155]. Some sources, notably [4, Subsections 9.1.1 and 10.1.1 together], also show the complete derivation of the (Taylor approximation of the) result, but the derivations are typically written specifically for the 2-D test problem. Another difficulty for casual reading of the derivations in some sources is also that they give - for didactic purposes - an interleaved presentation of the general theory and its application to the model problem. The key problem with the above is that it is not clear which results are valid in more generality than for the model problem and - even for the model problem - which part of the results extend to other than two dimensions.

Thus, this note has the dual purposes of reviewing the full derivation of the exact formula (1.3), while clearly separating the general theory from its application to the model problem, and extending the proof to the general $d$-D model problem for all $d=1,2,3, \ldots$. The final result is collected in Theorem 11 that confirms that (1.3) is indeed the optimal relaxation parameter for the SOR method in any space dimension.

This paper is organized as follows. Section 2 reviews the relevant results on the convergence theory of the SOR method that are valid for linear systems with general system matrices. Moreover, if the system matrix satisfies a certain technical assumption, the relation between spectral radius of its Jacobi and SOR iteration matrices can be established. Section 3 introduces the Poisson equation as model problem and its discretization by the finite difference method that results in the particular system matrices, for which we wish to derive the optimal relaxation parameter. To prepare this, the eigenvalue and eigenvectors of the system matrices are derived, first in 1-D and in 2-D, then furthermore in general $d$-D. For interest, a way of discrete Fourier transforms is showed to derive the eigenvalues of the 2-D model problem, whose idea is from [12, Section 4.6] and [2]. Section 4 shows that the model problem satisfies the needed technical assumption in Section 2 and then derives the optimal relaxation parameter for any $d$-D model problem, using the eigenvalues of the system matrices of the linear system from Section 3.

## 2 Some Results for the SOR Method for General System Matrices

For iterative methods, various authors have developed a convergence theory, see, e.g., [5], [11], [12], and others. Here, we state and derive some well-known theorems about the convergence of classical iterative methods. These statements apply to linear systems $A x=b$ with general $n \times n$ system matrices without making any assumption on $A$, aside from $A$ being non-singular. Then, under only one technical assumption, it is possible to derive a formula for the optimal relaxation parameter for the SOR method, if one has the spectral radius of the Jacobi method available.

Suppose that a matrix norm $\|\cdot\|$ is defined. Consider the iteration in the form of

$$
\begin{equation*}
M x^{(k+1)}=N x^{(k)}+b, \quad k=0,1,2, \ldots, \tag{2.1}
\end{equation*}
$$

where $A=M-N$ is a splitting of $A$ with a non-singular matrix $M$. Let the true solution be denoted as $x$, and the approximation solution in the $k$ th iteration be $x^{(k)}$. One obvious way to estimate the
error in each step is by analyzing $e^{(k)}=x-x^{(k)}$. From $A x=b$ with $A=M-N$ and from (2.1), we have

$$
\begin{aligned}
x & =M^{-1} N x+M^{-1} b, \\
x^{(k+1)} & =M^{-1} N x^{(k)}+M^{-1} b .
\end{aligned}
$$

Subtracting the two equations, we obtain $x-x^{(k+1)}=M^{-1} N\left(x-x^{(k)}\right)$, namely $e^{(k+1)}=M^{-1} N e^{(k)}=$ $G e^{(k)}$, with iteration matrix $G=M^{-1} N$. Substituting $e^{(k)}$ with its previous item inductively, we have $e^{(k)}=G^{k} e^{(0)}$, where $e^{(0)}$ is the initial error.

Theorem 1. [11, Section 7.3] An iterative method converges for any initial guess $x^{(0)}$, if and only if the spectral radius $\rho(G)$ of its iteration matrix $G$ satisfies

$$
\rho(G)<1,
$$

where $\rho(G)$ is the spectral radius of $G$. For large $k$, it holds that

$$
\frac{\left\|e^{(k+1)}\right\|}{\left\|e^{(k)}\right\|} \approx \rho(G) .
$$

Therefore, $\rho(G)$ is the convergence rate of the iterative method.
Here, the spectral radius $\rho(G)$ is the absolute value of the largest (in magnitude) eigenvalue of $G$, that is, $\rho(G)=\max _{k}\left|\lambda_{k}\right|$, if $\lambda_{k}, k=1, \ldots, n$, are the eigenvalues of $G$.

Consider now the classical iterative method SOR, which is given by (2.1) with the splitting matrix $M=\frac{1}{\omega} D-L$, where we use the standard notation of $A=D-L-U$ with $D$ denoting the diagonal part of $A$ and $-L$ and $-U$ denoting the strictly lower and upper triangular parts of $A$, respectively. The parameter $\omega$ is called a relaxation parameter. Let $G_{J}$ and $G_{\omega}$ denote the iteration matrices of the Jacobi and SOR methods, respectively. Here, $G_{J}=I-D^{-1} A=D^{-1}(L+U)$, and $G_{\omega}=$ $I-\left(\omega^{-1} D-L\right)^{-1} A=(D-\omega L)^{-1}[(1-\omega) D+\omega U]$.

Theorem 2. [5, Section 4.4] If the SOR method converges, then the parameter $\omega$ satisfies that $0<$ $\omega<2$.

Proof. The iteration matrix for the SOR method is given by

$$
G_{\omega}=M^{-1} N=(D-\omega L)^{-1}[(1-\omega) D+\omega U] .
$$

Using properties of the determinant and the fact that the matrices $(1-\omega) I+\omega D^{-1} U,\left(I-\omega D^{-1} L\right)$, and hence also $\left(I-\omega D^{-1} L\right)^{-1}$ are triangular, we can compute

$$
\begin{aligned}
\operatorname{det}\left(G_{\omega}\right) & =\operatorname{det}\left[(D-\omega L)^{-1}[(1-\omega) D+\omega U]\right] \\
& =\operatorname{det}\left[\left[\left(D\left(I-\omega D^{-1} L\right)\right]^{-1}\right] \operatorname{det}\left[D\left[(1-\omega) I+\omega D^{-1} U\right]\right]\right. \\
& =\operatorname{det}\left(D^{-1}\right) \operatorname{det}\left[\left(I-\omega D^{-1} L\right)^{-1}\right] \operatorname{det}(D) \operatorname{det}\left[(1-\omega) I+\omega D^{-1} U\right] \\
& =\operatorname{det}\left[(1-\omega) I+\omega D^{-1} U\right] \\
& =(1-\omega)^{n},
\end{aligned}
$$

where $n$ denotes the dimension of the matrix $A$. Therefore, the eigenvalues $\lambda_{i}, i=1, \ldots, n$, of the iteration matrix $G_{\omega}$ satisfy

$$
\begin{equation*}
\prod_{i=1}^{n} \lambda_{i}=(1-\omega)^{n} \tag{2.2}
\end{equation*}
$$

Hence, for the maximum eigenvalue in absolute value $\rho\left(G_{\omega}\right)$, it must be not less than $|1-\omega|$, otherwise (2.2) cannot hold. If the SOR is convergent, then $\rho\left(G_{\omega}\right)<1$. Hence, it follows that $|1-\omega|<1$. For real numbers $\omega$, this means that $0<\omega<2$.

Theorem 3. [4, Section 10.1] Let $A=D-L-U$ be a matrix that satisfies the technical assumption

$$
\begin{equation*}
\operatorname{det}\left(k D-\gamma L-\gamma^{-1} U\right)=\operatorname{det}(k D-L-U) \quad \text { for all } k, \gamma \in \mathbb{R} \backslash\{0\} \text {, } \tag{2.3}
\end{equation*}
$$

and $G_{J}$ and $G_{\omega}$ the iteration matrices of the Jacobi and SOR methods, as defined above. If $\mu$ is an eigenvalue of $G_{J}$ and $\lambda \neq 0$ satisfies

$$
\begin{equation*}
\mu=\frac{\lambda+\omega-1}{\omega \lambda^{1 / 2}} \tag{2.4}
\end{equation*}
$$

for some $\omega \in(0,2)$, then $\lambda$ is an eigenvalue of $G_{\omega}$.
Proof. We have the definitions for $G_{J}$ and $G_{\omega}$ that

$$
\begin{aligned}
& G_{J}=D^{-1}(L+U) \\
& G_{\omega}=(D-\omega L)^{-1}[(1-\omega) D+\omega U] .
\end{aligned}
$$

Let $\lambda$ be one eigenvalue of $G_{\omega}, \lambda \neq 0$. Because $U$ and $L$ are strictly upper and lower triangular matrices, and their main diagonal are zeros, we have $\operatorname{det}\left(D^{-1}\right) \operatorname{det}(D-\omega L)=1$.

$$
\begin{aligned}
\operatorname{det}\left(G_{\omega}-\lambda I\right) & =\operatorname{det}\left[(D-\omega L)^{-1}[(1-\omega) D+\omega U]-\lambda I\right] \\
& =\operatorname{det}\left(D^{-1}\right) \operatorname{det}(D-\omega L) \operatorname{det}\left[(D-\omega L)^{-1}[(1-\omega) D+\omega U]-\lambda I\right] \\
& =\operatorname{det}\left[(1-\omega) I+\omega D^{-1} U-\lambda I+\omega \lambda D^{-1} L\right] \\
& =\omega^{n} \operatorname{det}\left[\left(\frac{1}{\omega}-1\right) I+D^{-1} U-\frac{\lambda I}{\omega}+\lambda D^{-1} L\right] \\
& =(-1)^{n} \omega^{n} \lambda^{1 / 2} \operatorname{det}\left(D^{-1}\right) \operatorname{det}\left[\frac{\omega+\lambda-1}{\omega \lambda^{1 / 2}} D-\lambda^{-1 / 2} U-\lambda^{1 / 2} L\right], \quad \text { use (2.3) } \\
& =(-1)^{n} \omega^{n} \lambda^{1 / 2} \operatorname{det}\left[\frac{\omega+\lambda-1}{\omega \lambda^{1 / 2}} I-D^{-1} U-D^{-1} L\right] \\
& =\omega^{n} \lambda^{1 / 2} \operatorname{det}\left[G_{J}-\frac{\omega+\lambda-1}{\omega \lambda^{1 / 2}} I\right] .
\end{aligned}
$$

Because $\lambda \neq 0$ satisfies

$$
\mu=\frac{\lambda+\omega-1}{\omega \lambda^{1 / 2}},
$$

$\mu$ is an eigenvalue of $G_{J}$. In turn, when $\mu$ is an eigenvalue of $G_{J}$, and $\lambda$ satisfies that relation, then $\lambda$ is an eigenvalue of $G_{\omega}$.

Theorem 4. [4, Section 10.1] Assume $A=D-L-U$ is a matrix that satisfies (2.3), and assume that $G_{J}=I-D^{-1} A$ has only real eigenvalues and that $\beta \equiv \rho\left(G_{J}\right)<1$. Then the SOR iteration converges for every $\omega \in(0,2)$, and the spectral radius of the SOR matrix is

$$
\rho\left(G_{\omega}\right)= \begin{cases}\frac{1}{4}\left[\omega \beta+\sqrt{(\omega \beta)^{2}-4(\omega-1)}\right]^{2} & \text { for } 0<\omega \leq \omega_{\text {opt }},  \tag{2.5}\\ w-1 & \text { for } \omega_{\text {opt }} \leq \omega<2,\end{cases}
$$

where $\omega_{\text {opt }}$, the optimal value of $\omega$, is

$$
\begin{equation*}
\omega_{o p t}=\frac{2}{1+\sqrt{1-\beta^{2}}} . \tag{2.6}
\end{equation*}
$$

For any other value of $\omega$, we have

$$
\begin{equation*}
\rho\left(G_{\omega_{\text {opt }}}\right)<\rho\left(G_{\omega}\right), \quad \text { for } \omega \in(0,2) \backslash\left\{\omega_{\text {opt }}\right\} . \tag{2.7}
\end{equation*}
$$

Proof. For a given $\omega, \rho\left(G_{\omega}\right)$ is the largest eigenvalue of $G_{\omega}$ in absolute value. Suppose that $\mu$ is an eigenvalue of $G_{J}$. Solving (2.4) in Theorem 3, we have

$$
\begin{equation*}
\lambda=\frac{1}{4}\left(\omega \mu \pm \sqrt{(\omega \mu)^{2}-4(\omega-1)}\right)^{2} . \tag{2.8}
\end{equation*}
$$

Then by Theorem 3, (2.8) gives two eigenvalues for $G_{\omega}$.
First, if $(\omega \mu)^{2}-4(\omega-1)<0$, then $\lambda$ is imaginary, and the absolute value of $\lambda$ is

$$
|\lambda|=\left(\sqrt{\frac{1}{4} \omega^{2} \mu^{2}+\omega-1-\frac{1}{4} \omega^{2} \mu^{2}}\right)^{2}=\omega-1,
$$

when

$$
\tilde{\omega} \equiv \frac{2\left(1-\sqrt{1-\mu^{2}}\right)}{\mu^{2}}<\omega<2 .
$$

Here, $\rho\left(G_{\omega}\right)=|\lambda|$ is independent of $\mu=\rho\left(G_{J}\right)$.
Second, if $(\omega \mu)^{2}-4(\omega-1) \geq 0$, then

$$
\rho\left(G_{\omega}\right)=\max _{\mu \in \sigma\left(G_{J}\right)} \frac{1}{4}\left[\omega|\mu|+\sqrt{(\omega|\mu|)^{2}-4(\omega-1)}\right]^{2}, \quad \omega \in(0, \tilde{\omega}] .
$$

On one hand, for a fixed $\omega, \rho\left(G_{\omega}\right)$ is an increasing function with respect to variable $|\mu|$. Hence, to get the spectral radius of $G_{\omega}$, let $|\mu|=\rho\left(G_{J}\right)=\beta$, we have

$$
\rho\left(G_{\omega}\right)=\frac{1}{4}\left[\omega \beta+\sqrt{(\omega \beta)^{2}-4(\omega-1)}\right]^{2} .
$$

On the other hand, $\rho\left(G_{\omega}\right)$ can be proved to be a decreasing function with respect of $\omega$ in $(0, \tilde{\omega}]$. Details are as follows. We have the first order derivative of $\rho\left(G_{\omega}\right)$

$$
\rho^{\prime}\left(G_{\omega}\right)=\frac{1}{2}\left(\omega \beta+\sqrt{(\omega \beta)^{2}-4(\omega-1)}\right)\left(\beta+\frac{\beta^{2} \omega-2}{\sqrt{(\omega \beta)^{2}-4(\omega-1)}}\right) .
$$

To determine its sign, we need to examine the sign of

$$
\beta+\frac{\beta^{2} \omega-2}{\sqrt{(\omega \beta)^{2}-4(\omega-1)}} .
$$

Because $\sqrt{(\omega \beta)^{2}-4(\omega-1)}>0, \beta<1, \omega<2$ and

$$
\begin{aligned}
\beta \sqrt{(\omega \beta)^{2}-4(\omega-1)}+\beta^{2} \omega-2 & <\sqrt{\omega^{2}-4 \omega+4}+\omega-2 \\
& =\sqrt{(\omega-2)^{2}}+\omega-2 \\
& =2-\omega+\omega-2 \\
& =0 .
\end{aligned}
$$

Therefore, $\rho^{\prime}\left(G_{\omega}\right)<0$ in the interval $(0, \tilde{\omega}]$, which implies that $\rho\left(G_{\omega}\right)$ is a decreasing function of $\omega$. When $\omega=\tilde{\omega}, \rho\left(G_{\omega}\right)$ gets its minimum in the interval $(0, \tilde{\omega}]$. We have proved above that in the interval $(\tilde{\omega}, 2), \rho\left(G_{\omega}\right)=\omega-1$ is an increasing function and also gets its minimum when $\omega$ approaches to $\tilde{\omega}$. Moreover, $\rho\left(G_{\omega}\right)$ is continuous at the point $\tilde{\omega}$. Considering that the optimal parameter $\omega$ is the very number that makes $\rho\left(G_{\omega}\right)$ gets its minimum. Therefore, $\tilde{\omega}=\omega_{\text {opt }}=\frac{2\left(1-\sqrt{1-\beta^{2}}\right)}{\beta^{2}}=\frac{2}{1+\sqrt{1-\beta^{2}}}$. According to the statement above, for any value of $\omega$,

$$
\rho\left(G_{\omega_{o p t}}\right)<\rho\left(G_{\omega}\right), \quad \omega \in(0,2) \backslash\left\{\omega_{\text {opt }}\right\} .
$$



Figure 1: Plot of $\rho\left(G_{\omega}\right)$ defined in (2.5).

To get a visual impression of the influence of the choice of $\omega$ on the spectral radius $\rho\left(G_{\omega}\right)$ given by the function defined in (2.5), we plot $\rho\left(G_{\omega}\right)$ as a function of $\omega$ in Figure 1 for several values of $\beta=\rho\left(G_{J}\right)$. The plot shows that the curve has a vertical asymptote from the left at $\omega=\omega_{\text {opt }}$, but is linear for $\omega>\omega_{\text {opt }}$. Thus, it is better to slightly overestimate $\omega_{\text {opt }}$ than to underestimate it, which is a classical rule-of-thumb for the SOR method [1, page 294].

## 3 The Model Problem and its Eigenvalues

### 3.1 The 1-D Model Problem

Many large-scale linear equation systems arise from the discretization of partial differential equations. One fundamental partial differential equation is the Poisson equation (1.1) subject to the boundary condition (1.2). In the model problem, we consider $g=0$ for simplicity and the unit interval $\Omega=(0,1)$ for 1-D case and the unit square $\Omega=(0,1) \times(0,1)$ for 2 -D case.

The finite difference method is now used to discretize this model, that is, to approximate the derivatives in the differential equation. Topics on discretization of PDEs can be found on many books, e.g., [7], [9], and [10]. First look at the Taylor series,

$$
\begin{align*}
& u(x+h)=u(x)+h \frac{d u}{d x}+\frac{h^{2}}{2!} \frac{d^{2} u}{d x^{2}}+\frac{h^{3}}{3!} \frac{d^{3} u}{d x^{3}}+\mathcal{O}\left(h^{4}\right),  \tag{3.1}\\
& u(x-h)=u(x)-h \frac{d u}{d x}+\frac{h^{2}}{2!} \frac{d^{2} u}{d x^{2}}-\frac{h^{3}}{3!} \frac{d^{3} u}{d x^{3}}+\mathcal{O}\left(h^{4}\right) . \tag{3.2}
\end{align*}
$$

Combining (3.1) and (3.2), we have

$$
\begin{equation*}
\frac{d^{2} u}{d x^{2}}=\frac{u(x-h)-2 u(x)+u(x+h)}{h^{2}}+\mathcal{O}\left(h^{2}\right) . \tag{3.3}
\end{equation*}
$$

This formula involves the solution at three points $u(x-h), u(x), u(x+h)$, hence this finite difference is said to use a three-point stencil. It is called a centered difference approximation of the second derivative. If any of the points is on the boundary, its value is known, namely as 0 in the model problem, otherwise the points are in the interior of $\Omega$.

Now, use the three-point stencil for the finite difference approximation

$$
\begin{equation*}
-u^{\prime \prime}=f(x) \quad \text { in } \Omega=(0,1) \tag{3.4}
\end{equation*}
$$

subject to the boundary condition

$$
\begin{equation*}
u(0)=0 \text { and } u(1)=0 . \tag{3.5}
\end{equation*}
$$

Let $h=1 /(N+1)$ denote the mesh spacing, where $N$ means dividing the $[0,1]$ interval with $N+2$ points. Hence, we get a mesh of uniformly spaced points $x_{k}=k h$ on the interval $[0,1]$, for $k=0,1, \ldots, N, N+1$. Then we can approximate $u^{\prime \prime}\left(x_{k}\right)$ at the interior points of the mesh with the formula according to (3.3)

$$
\begin{equation*}
h^{2} f\left(x_{k}\right)=-h^{2} u^{\prime \prime}\left(x_{k}\right) \approx-u\left(x_{k-1}\right)+2 u\left(x_{k}\right)-u\left(x_{k+1}\right), \quad k=1, \ldots, N, \tag{3.6}
\end{equation*}
$$

where we use the boundary conditions $u_{0}=u_{N+1}=0$. Therefore, we get a system of $N$ equations in $N$ unknowns.

If the unknowns $u_{k} \approx u\left(x_{k}\right)$ are organized in a column vector $u=\left[u_{1}, u_{1}, \ldots, u_{N}\right]^{T} \in \mathbb{R}^{N}$, the equations in (3.6) can be organized in matrix form

$$
A_{N} u=b
$$

with the tri-diagonal system matrix

$$
A_{N}=\left[\begin{array}{cccccc}
2 & -1 & & & &  \tag{3.7}\\
-1 & 2 & -1 & & & \\
& -1 & 2 & -1 & & \\
& & \ddots & \ddots & \ddots & \\
& & & -1 & 2 & -1 \\
& & & & -1 & 2
\end{array}\right]
$$

and right-hand side $b=h^{2}\left[f\left(x_{1}\right), f\left(x_{2}\right), \ldots, f\left(x_{N}\right)\right]^{T}$.
Before moving on to the proof of the following theorem that gives the eigenvalues and eigenvectors of the system matrix in (3.7), we quote one result for the solution to finite difference equations for simple roots of the characteristic polynomial from [7, Section 1.3].

Theorem 5. If all the roots of the characteristic polynomial of a linear difference equation are simple and nonzero, then each solution of the difference equation is a linear combination of such special solution.

In other words, if $r_{1}, r_{2}, \ldots, r_{m}$ are simple and nonzero roots for characteristic polynomial of a difference equation, then the components $x_{i}$ of the solution vector of the difference equation $x$ can be written as

$$
x_{i}=\sum_{i=1}^{m} a_{k} r_{k}^{i}, \quad 1 \leq i \leq m
$$

where $a_{k}$ are coefficients that are determined from the boundary conditions.
Theorem 6. Let $A_{N} \in \mathbb{R}^{N \times N}$ be the matrix given in (3.7), then the eigenvalues $\lambda_{k}$ and eigenvectors $z^{(k)}, k=1, \ldots, N$, of $A_{N}$ are given by

$$
\begin{equation*}
\lambda_{k}=2\left(1-\cos \left(\frac{k \pi}{N+1}\right)\right)=4 \sin ^{2}\left(\frac{k \pi}{2(N+1)}\right) \tag{3.8}
\end{equation*}
$$

and the components of $z^{(k)}=\left(z_{\ell}^{(k)}\right)$ by

$$
\begin{equation*}
z_{\ell}^{(k)}=\sqrt{\frac{2}{N+1}} \sin \left(\frac{k \ell \pi}{N+1}\right), \quad \ell=1,2, \ldots, N \tag{3.9}
\end{equation*}
$$

Proof. Let $(\lambda, z)$ denote one eigenpair. So, $A_{N} z=\lambda z$ or $\left(A_{N}-\lambda I\right) z=0$, namely

$$
\left[\begin{array}{ccccc}
2-\lambda & -1 & & & \\
-1 & 2-\lambda & -1 & & \\
& \ddots & \ddots & \ddots & \\
& & -1 & 2-\lambda & -1 \\
& & & -1 & 2-\lambda
\end{array}\right]\left[\begin{array}{c}
z_{1} \\
z_{2} \\
z_{3} \\
\vdots \\
z_{N}
\end{array}\right]=0
$$

The $k$ th component of this equation is $-z_{k-1}+(2-\lambda) z_{k}-z_{k+1}=0, k=1, \ldots, N$, with the boundary conditions $z_{0}=z_{N+1}=0$. We can re-write this as

$$
\begin{equation*}
z_{k+2}-(2-\lambda) z_{k+1}+z_{k}=0, \tag{3.10}
\end{equation*}
$$

where $k=0, \ldots, N-1$, and $z_{0}=z_{N+1}=0$. This problem is a second-order linear finite difference equation with constant coefficients, one of whom contains the eigenvalue $\lambda$ that also needs to be determined as part of the problem. We have the characteristic polynomial $p(r)=r^{2}-(2-\lambda) r+1$ of (3.10).

If the characteristic equation had a double root $r_{1}$, then the general solution of (3.10) had the form $z_{\ell}=\alpha r_{1}^{\ell}+\beta \ell r_{1}^{\ell}$. The boundary conditions $z_{0}=\alpha=0$ and $z_{N+1}=\alpha r_{1}^{N+1}+\beta(N+1) r_{1}^{N+1}=0$ would then give $\alpha=\beta=0$. Hence, we need to have two distinct roots $r_{1} \neq r_{2}$ and have the general solution $z_{\ell}=\alpha r_{1}^{\ell}+\beta r_{2}^{\ell}$, by Theorem 5 .

The boundary condition then read

$$
\begin{cases}0 & =\alpha+\beta, \\ 0 & =\alpha r_{1}^{N+1}+\beta r_{2}^{N+1} .\end{cases}
$$

Solving these equations, we have

$$
\left(\frac{r_{1}}{r_{2}}\right)^{N+1}=1
$$

that is the fraction $r_{1} / r_{2}$ can be any of the $N+1$ complex roots of unity given by

$$
\begin{equation*}
\frac{r_{1}}{r_{2}}=e^{\frac{2 \pi k i}{N+1}}, \quad k=0, \ldots, N \tag{3.11}
\end{equation*}
$$

or $r_{1}=r_{2} e^{\frac{2 \pi k i}{N+1}}$. Since $r_{1}$ and $r_{2}$ are the roots of the characteristic polynomial, we also have $p(r)=$ $r^{2}-(2-\lambda) r+1=\left(r-r_{1}\right)\left(r-r_{2}\right)=r^{2}-\left(r_{1}+r_{2}\right) r+r_{1} r_{2}$. Matching the constant terms gives $1=r_{1} r_{2}$ and thus

$$
r_{1}=e^{\frac{i k \pi}{N+1}}, \quad r_{2}=e^{\frac{-i k \pi}{N+1}}, \quad k=1, \ldots, N
$$

notice that the case $k=0$ would result in a double root, and we showed above already that this cannot be the case. Matching now the coefficient of the linear term gives $2-\lambda=r_{1}+r_{2}$, which determines the eigenvalues $\lambda_{k}, k=1, \ldots, N$, as

$$
\lambda_{k}=2-\left(r_{1}+r_{2}\right)=2-\left(e^{\frac{i k \pi}{N+1}}+e^{\frac{-i k \pi}{N+1}}\right)=2-2 \cos \left(\frac{k \pi}{N+1}\right), \quad k=1,2, \ldots, N,
$$

using the formula $\cos \theta=\left(e^{i \theta}+e^{-i \theta}\right) / 2$. Using moreover the trigonometric identity $\cos 2 x=1-2 \sin ^{2} x$ gives the alternative form $\lambda_{k}=4 \sin ^{2}\left(\frac{k \pi}{2(N+1)}\right)$.

Therefore, using $\sin \theta=\left(e^{i \theta}-e^{-i \theta}\right) /(2 i)$, we finally have the components of the eigenvectors

$$
z_{\ell}^{(k)}=\alpha r_{1}^{\ell}+\beta r_{2}^{\ell}=\alpha\left(e^{\frac{k \ell \pi i}{N+1}}-e^{\frac{-k \ell \pi i}{N+1}}\right)=2 i \alpha \sin \left(\frac{k \ell \pi}{N+1}\right)
$$

Letting $\alpha=\sqrt{\frac{2}{N+1}} /(2 i)$, we have

$$
z_{\ell}^{(k)}=\sqrt{\frac{2}{N+1}} \sin \left(\frac{k \ell \pi}{N+1}\right)
$$

The choice of $\alpha$ in the scaling of the eigenvectors ensures that they are real and are normalized, that is, $\left\|z^{(k)}\right\|_{2}=1$. Notice moreover that the eigenvectors form an orthogonal set, since $A_{N}$ is symmetric. Together with the normalization, this gives an orthonormal set.

We will later need a generalization to tri-diagonal matrices of the following form.

Theorem 7. [4, Section 9.1] Let $A$ be an $N \times N$ matrix of the form

$$
\left[\begin{array}{ccccc}
\alpha & \beta & & & \\
\beta & \alpha & \beta & & \\
& \ddots & \ddots & \ddots & \\
& & \beta & \alpha & \beta \\
& & & \beta & \alpha
\end{array}\right]
$$

Then the eigenvalues $\lambda_{k}$ and eigenvectors $z^{(k)}, k=1, \ldots, N$, of $A$ are given by

$$
\begin{equation*}
\lambda_{k}=\alpha+2 \beta \cos \left(\frac{k \pi}{N+1}\right) \tag{3.12}
\end{equation*}
$$

and the m-th component of $z^{(k)}=\left(z_{m}^{(k)}\right)$ by

$$
\begin{equation*}
z_{m}^{(k)}=\sqrt{\frac{2}{N+1}} \sin \left(\frac{m k \pi}{N+1}\right), \quad k, m=1,2, \ldots, N \tag{3.13}
\end{equation*}
$$

Proof. See the proof of Lemma 9.1.1 in Greenbaum [4, Section 9.1].

Notice that the eigenvectors do not involve the coefficients $\alpha$ and $\beta$, thus all matrices $A$ of this form can be diagonalized using the same transformation matrix $Q$ with the eigenvectors $z^{(k)}$ as its columns. Since the eigenvectors form an orthonormal set, this matrix will be orthogonal, that is, it satisfies $Q^{T} Q=I$ or $Q^{-1}=Q^{T}$.

### 3.2 The 2-D Model Problem

In two dimensions, let the domain of the model problem be the unit square $\Omega=(0,1) \times(0,1)$. The Poisson problem (1.1)-(1.2) can be written in 2-D concretely as

$$
\begin{equation*}
-\frac{\partial^{2} u}{\partial x^{2}}-\frac{\partial^{2} u}{\partial y^{2}}=f \quad \text { in } \Omega \tag{3.14}
\end{equation*}
$$

subject to the boundary condition

$$
\begin{equation*}
u=0 \quad \text { on } \partial \Omega \tag{3.15}
\end{equation*}
$$

Let $h=1 /(N+1)$ denote the mesh spacing, where $N$ means divides each direction of $[0,1] \times[0,1]$ into $N+2$ points. Hence, we get a mesh of uniformly spaced points $\left(x_{k}, y_{\ell}\right)=(k h, \ell h)$ on the bounded region $\bar{\Omega}$, for $k, \ell=0,1, \ldots, N, N+1$. According to (3.3), fixing the $y$ variable, we have

$$
\begin{equation*}
\frac{\partial^{2} u}{\partial x^{2}}\left(x_{k}, y_{\ell}\right) \approx \frac{u\left(x_{k-1}, y_{\ell}\right)-2 u\left(x_{k}, y_{\ell}\right)+u\left(x_{k+1}, y_{\ell}\right)}{h^{2}} \tag{3.16}
\end{equation*}
$$

With the $x$ variable being fixed, we obtain

$$
\begin{equation*}
\frac{\partial^{2} u}{\partial y^{2}}\left(x_{k}, y_{\ell}\right) \approx \frac{u\left(x_{k}, y_{\ell-1}\right)-2 u\left(x_{k}, y_{\ell}\right)+u\left(x_{k}, y_{\ell+1}\right)}{h^{2}} . \tag{3.17}
\end{equation*}
$$

Combining (3.16) and (3.17) together, we can approximate $-\triangle u\left(x_{k}, y_{\ell}\right)=f\left(x_{k}, y_{\ell}\right)$ with the following formula

$$
\begin{equation*}
-u_{k-1, \ell}-u_{k, \ell-1}+4 u_{k, \ell}-u_{k, \ell+1}-u_{k+1, \ell}=h^{2} f_{k, \ell}, \quad k, \ell=1,2, \ldots, N . \tag{3.18}
\end{equation*}
$$

Using the fact that the value of $u=0$ on the boundary is known and we only have to consider approximations $u_{k \ell} \approx u\left(x_{k}, y_{\ell}\right)$ for the interior points, we have $N^{2}$ equations in $N^{2}$ unknowns. If the unknowns are organized in natural ordering into the column vector

$$
u=\left[u_{11}, u_{21}, \ldots, u_{N 1}, u_{12}, u_{22}, \ldots, u_{N 2}, \ldots, u_{1 N}, u_{2 N}, \ldots, u_{N N}\right]^{T} \in \mathbb{R}^{N^{2}}
$$

then the problem can be stated in matrix form

$$
\begin{equation*}
A u=b \tag{3.19}
\end{equation*}
$$

with system matrix

$$
A=\left[\begin{array}{ccccc}
S & T & & &  \tag{3.20}\\
T & S & T & & \\
& \ddots & \ddots & \ddots & \\
& & T & S & T \\
& & & T & S
\end{array}\right], S=\left[\begin{array}{ccccc}
4 & -1 & & & \\
-1 & 4 & -1 & & \\
& \ddots & \ddots & \ddots & \\
& & -1 & 4 & -1 \\
& & & -1 & 4
\end{array}\right], T=\left[\begin{array}{lllll}
-1 & & & & \\
& -1 & & & \\
& & \ddots & & \\
& & & -1 & \\
& & & & -1
\end{array}\right]
$$

and $b$ collecting the terms $h^{2} f\left(x_{k}, y_{\ell}\right)$ in the same order as the terms in $u$.
Theorem 8. [4, Section 9.1] Let $A \in \mathbb{R}^{N^{2} \times N^{2}}$ be defined as in (3.20). The eigenvalues of $A$ are then

$$
\lambda_{k, \ell}=4 \sin ^{2}\left(\frac{k \pi}{2(N+1)}\right)+4 \sin ^{2}\left(\frac{\ell \pi}{2(N+1)}\right), \quad k, \ell=1,2, \ldots, N,
$$

with corresponding eigenvectors, with components counted with respect to the mesh point ( $m, n$ ),

$$
z_{m, n}^{(k, \ell)}=\frac{2}{N+1} \sin \left(\frac{m k \pi}{N+1}\right) \sin \left(\frac{n \ell \pi}{N+1}\right), \quad m, n, k, \ell=1,2, \ldots, N
$$

Proof. Let $(\lambda, z)$ be one eigenpair of $A$, namely $A z=\lambda z . A$ is in block form, having $N \times N$ blocks of size $N \times N$. Let eigenvector $z$ be partitioned as

$$
z=\left(\begin{array}{c}
z_{1} \\
\vdots \\
z_{N}
\end{array}\right) \in \mathbb{R}^{N^{2}} \quad \text { with the } n \text {-th block of } z \text { being } \quad z_{n}=\left(\begin{array}{c}
z_{1, n} \\
\vdots \\
z_{N, n}
\end{array}\right) \in \mathbb{R}^{N}, \quad n=1, \ldots, N .
$$

According to the block matrix eigenproblem, we have

$$
\begin{equation*}
T z_{n-1}+(S-\lambda I) z_{n}+T z_{n+1}=0, \quad n=1,2, \ldots, N \tag{3.21}
\end{equation*}
$$

where $z_{0}=z_{N+1}=0$.
In Section 3.1, we have seen that $S$ and $T$ can be diagonalized as $S=Q \Lambda_{S} Q^{T}$ and $T=Q \Lambda_{T} Q^{T}$ with the same orthogonal transformation matrix $Q$ and diagonal matrices $\Lambda_{S}$ and $\Lambda_{T}$, whose $k$-th diagonal entries are

$$
\begin{aligned}
& \lambda_{S, k}=4-2 \cos \left(\frac{k \pi}{N+1}\right), \\
& \lambda_{T, k}=-1 .
\end{aligned}
$$

By (3.13), the $m$ th component of column $k$ of $Q$ is

$$
q_{m}^{(k)}=\sqrt{\frac{2}{N+1}} \sin \left(\frac{m k \pi}{N+1}\right), \quad m, k=1,2, \ldots, N
$$

Using (3.21), we have

$$
\begin{aligned}
& Q^{T} T z_{n-1}+Q^{T}(S-\lambda I) z_{n}+Q^{T} T z_{n+1}=0 \\
\Rightarrow & \Lambda_{T} y_{n-1}+\left(\Lambda_{S}-\lambda I\right) y_{n}+\Lambda_{T} y_{n+1}=0, \quad y_{n}=Q^{T} z_{n}, \quad n=1,2, \ldots, N .
\end{aligned}
$$

Both $\Lambda_{T}$ and $\left(\Lambda_{S}-\lambda I\right)$ are diagonal. So we have the following difference equation

$$
\lambda_{T, k} y_{k, n-1}+\lambda_{S, k} y_{k, n}+\lambda_{T, k} y_{k, n+1}=\lambda y_{k, n}, \quad k=1,2, \ldots, N .
$$

This problem can be rewritten in matrix form, when the index $k$ is fixed and all components of $y$ except the $k$ th one are 0 ,

$$
\left(\begin{array}{ccccc}
\lambda_{S, k} & \lambda_{T, k} & & & \\
\lambda_{T, k} & \lambda_{S, k} & \lambda_{T, k} & & \\
& \ddots & \ddots & \ddots & \\
& & \lambda_{T, k} & \lambda_{S, k} & \lambda_{T, k} \\
& & & \lambda_{T, k} & \lambda_{S, k}
\end{array}\right)\left(\begin{array}{c}
y_{k, 1} \\
y_{k, 2} \\
\vdots \\
y_{k, N-1} \\
y_{k, N}
\end{array}\right)=\lambda\left(\begin{array}{c}
y_{k, 1} \\
y_{k, 2} \\
\vdots \\
y_{k, N-1} \\
y_{k, N}
\end{array}\right)
$$

By Theorem 7, the eigenpairs for such an eigenproblem are given as

$$
\begin{align*}
\lambda_{k, \ell} & =\lambda_{S, k}+2 \lambda_{T, k} \cos \left(\frac{\ell \pi}{N+1}\right) \\
& =4-2 \cos \left(\frac{k \pi}{N+1}\right)-2 \cos \left(\frac{\ell \pi}{N+1}\right)  \tag{3.22}\\
& =4 \sin ^{2}\left(\frac{k \pi}{2(N+1)}\right)+4 \sin ^{2}\left(\frac{\ell \pi}{2(N+1)}\right),  \tag{3.23}\\
y_{k, n}^{(k, \ell)} & =\sqrt{\frac{2}{N+1}} \sin \left(\frac{n \ell \pi}{N+1}\right),
\end{align*}
$$

where $k, \ell, n=1,2, \ldots, N$. Here, $y_{k, n}^{(k, \ell)}$ is the $k$ th entry of the $n$th block of the corresponding eigenvector.
Since $y_{n}=Q^{T} z_{n}, n=1,2, \ldots, N$, and because all components of $y$ except the $k$ th one are 0 , only the $k$ th entry of the $n$th block of $y$ is nonzero. Hence, we have

$$
z_{m, n}^{(k, \ell)}=q_{m}^{k} y_{k, n}^{(k, \ell)}=\frac{2}{N+1} \sin \left(\frac{m k \pi}{N+1}\right) \sin \left(\frac{n \ell \pi}{N+1}\right) .
$$

Remark. There is another way to find the eigenvalues of the 2-D model problem. For interest, we show in detail how to use discrete Fourier transforms to get the eigenvalues. The idea comes from $[12$, Section 4.6] and [2]. For more details about Fourier transformations, [8] is a good reference.

Let $A$ be the coefficient matrix given by (3.20). To determine the eigenvalues of $A$, we can solve the equation

$$
\begin{equation*}
-u_{k-1, \ell}-u_{k, \ell-1}+(4-\lambda) u_{k, \ell}-u_{k, \ell+1}-u_{k+1, \ell}=0 \tag{3.24}
\end{equation*}
$$

where $u_{0, \ell}=0, u_{k, 0}=0$, for $k, \ell=1,2, \ldots, N$. By [8, Section 19.4], the 2-D inverse sine transform for $u_{k, \ell}$ is

$$
u_{k, \ell}=\frac{4}{(N+1)^{2}} \sum_{m=1}^{N} \sum_{n=1}^{N} \hat{u}_{m, n} \sin \frac{k \pi m}{N+1} \sin \frac{\ell \pi n}{N+1},
$$

which satisfies the Dirichlet boundary condition $u=0$ on the boundary, when $k=0, N+1$ or $\ell=0, N+1$. Then we introduce some notations, following [2]. Denote

$$
\begin{aligned}
S_{b}^{a} & =\sin \frac{a \pi b}{N+1}, \\
C_{b}^{a} & =\cos \frac{a \pi b}{N+1} .
\end{aligned}
$$

Applying this to (3.24), we get

$$
\frac{4}{(N+1)^{2}} \sum_{m=1}^{N} \sum_{n=1}^{N} \hat{u}_{m, n}\left(S_{m}^{k-1} S_{n}^{\ell}+S_{m}^{k} S_{n}^{\ell-1}+S_{m}^{k} S_{n}^{\ell+1}+S_{m}^{k+1} S_{n}^{\ell}\right)=\frac{4(4-\lambda)}{(N+1)^{2}} \sum_{m=1}^{N} \sum_{n=1}^{N} \hat{u}_{m, n} S_{m}^{k} S_{n}^{\ell} .
$$

Canceling out the common factor and taking away the summation, we have

$$
\begin{equation*}
S_{m}^{k-1} S_{n}^{\ell}+S_{m}^{k} S_{n}^{\ell-1}+S_{m}^{k} S_{n}^{\ell+1}+S_{m}^{k+1} S_{n}^{\ell}=(4-\lambda) S_{m}^{k} S_{n}^{\ell} \tag{3.25}
\end{equation*}
$$

Since $\sin (a+b)=\sin a \cos b+\cos a \sin b$, we have in our notation

$$
S_{m}^{a+b}=S_{m}^{a} C_{m}^{b}+C_{m}^{a} S_{m}^{b}
$$

Moreover,

$$
\begin{aligned}
C_{m}^{-a} & =C_{m}^{a}, \\
S_{m}^{-a} & =-S_{m}^{a} .
\end{aligned}
$$

Therefore,

$$
\begin{aligned}
& S_{m}^{k+1} S_{n}^{\ell}+S_{m}^{k-1} S_{n}^{\ell}=S_{n}^{\ell}\left(S_{m}^{k} C_{m}^{1}+C_{m}^{k} S_{m}^{1}+S_{m}^{k} C_{m}^{1}-C_{m}^{k} S_{m}^{1}\right)=2 S_{n}^{\ell} S_{m}^{k} C_{m}^{1} \\
& S_{m}^{k} S_{n}^{\ell+1}+S_{m}^{k} S_{n}^{\ell-1}=S_{m}^{k}\left(S_{n}^{\ell} C_{n}^{1}+C_{n}^{\ell} S_{n}^{1}+S_{n}^{\ell} C_{n}^{1}-C_{n}^{\ell} S_{n}^{1}\right)=2 S_{m}^{k} S_{n}^{\ell} C_{n}^{1}
\end{aligned}
$$

Using this, (3.25) can be simplified to be

$$
\begin{aligned}
2 S_{n}^{\ell} S_{m}^{k} C_{m}^{1}+2 S_{m}^{k} S_{n}^{\ell} C_{n}^{1} & =(4-\lambda) S_{m}^{k} S_{n}^{\ell} \\
2 C_{m}^{1}+2 C_{n}^{1} & =4-\lambda .
\end{aligned}
$$

Therefore, we have

$$
\lambda_{m, n}=4-2 C_{m}^{1}-2 C_{n}^{1}, \quad m, n=1,2, \ldots, N .
$$

Here, $m$ and $n$ can be replaced with index $k$ and $\ell$, and the formula for the eigenvalues of the model problem in 2-D can be written as

$$
\begin{aligned}
\lambda_{k, \ell} & =4-2 \cos \left(\frac{k \pi}{N+1}\right)-2 \cos \left(\frac{\ell \pi}{N+1}\right) \\
& =4 \sin ^{2}\left(\frac{k \pi}{2(N+1)}\right)+4 \sin ^{2}\left(\frac{\ell \pi}{2(N+1)}\right), \quad k, \ell=1,2, \ldots, N .
\end{aligned}
$$

It can be compared with (3.23) and is the same as that result. The Fourier transformation also provides a way to solve the Poisson equation. Interested readers can refer to [8].

### 3.3 The $d$-D Model Problem

In some applications, we need to discretize the model problem in 3-D domain. Furthermore, we can generalize the result for $d$-D case of the Poisson equation (1.1) with boundary condition (1.2). In this section, we also consider $g=0$ and the unit region $\Omega=(0,1) \times(0,1) \times \ldots \times(0,1)$ for $d$-D case, which can be written in $d$-D case as

$$
\begin{equation*}
-\frac{\partial^{2} u}{\partial x_{1}^{2}}-\frac{\partial^{2} u}{\partial x_{2}^{2}}-\cdots-\frac{\partial^{2} u}{\partial x_{d}^{2}}=f \quad \text { in } \Omega \tag{3.26}
\end{equation*}
$$

subject to the boundary condition

$$
\begin{equation*}
u=0 \quad \text { on } \partial \Omega \tag{3.27}
\end{equation*}
$$

Similar to the previous cases, we can derive the system matrix $A$ for the $d$-D case, where $A u=b$. The system matrix

$$
\begin{gather*}
A_{N^{d} \times N^{d}}=\left[\begin{array}{ccccc}
S_{1} & T & & & \\
T & S_{1} & T & & \\
& \ddots & \ddots & \ddots & \\
& & T & S_{1} & T \\
& & & T & S_{1}
\end{array}\right], \quad S_{i}=\left[\begin{array}{ccccc}
S_{i+1} & T & & & \\
T & S_{i+1} & T & & \\
& \ddots & \ddots & \ddots & \\
& & T & S_{i+1} & T \\
& & & T & S_{i+1}
\end{array}\right],  \tag{3.28}\\
S_{d-1}=\left[\begin{array}{cccccc}
2 d & -1 & & & \\
-1 & 2 d & -1 & & \\
& \ddots & \ddots & \ddots & \\
& & -1 & 2 d & -1 \\
& & & -1 & 2 d
\end{array}\right], \quad T=\left[\begin{array}{ccccc}
-1 & & & & \\
& -1 & & & \\
& & \ddots & & \\
& & & -1 & \\
& & & & -1
\end{array}\right],
\end{gather*}
$$

where $1 \leq i \leq d-2$. Here, $A$ is of size $N^{d} \times N^{d}$ and $S_{d-1}$ of size $N \times N$. The matrices $T=-I$ are negative identity matrices with dimensions adjusted implicitly according to the size of $S_{i+1}$.

Theorem 9. Let $A \in \mathbb{R}^{N^{d} \times N^{d}}$ be defined as in (3.28). The eigenvalues $\lambda_{k_{1}, k_{2}, \ldots, k_{d}}$ and eigenvectors $z_{m_{1}, m_{2}, \ldots, m_{d}}^{\left(k_{1}, k_{2}, \ldots, k_{d}\right)}$ of $A$,, counted with respect to the mesh points $\left(k_{1}, k_{2}, \ldots, k_{d}\right)$, are given by

$$
\begin{equation*}
\lambda_{k_{1}, k_{2}, \ldots, k_{d}}=2 d-2 \sum_{i=1}^{d} \cos \left(\frac{k_{i} \pi}{N+1}\right)=4 \sum_{i=1}^{d} \sin ^{2}\left(\frac{k_{i} \pi}{2(N+1)}\right) \tag{3.29}
\end{equation*}
$$

and components of the eigenvectors, counted with respect to the mesh point $\left(m_{1}, m_{2}, \ldots, m_{d}\right)$,

$$
z_{m_{1}, m_{2}, \ldots, m_{d}}^{\left(k_{1}, k_{2}, \ldots, k_{d}\right)}=\left(\frac{2}{N+1}\right)^{\frac{d}{2}} \prod_{i=1}^{d} \sin \left(\frac{k_{i} m_{i} \pi}{N+1}\right), \quad k_{i}, m_{i}=1,2, \ldots, N .
$$

In (3.29), the term $2 d$ is determined by the main diagonal entry of $A$.
Proof. To prove it by induction, we have the initial step that $d=1$ is true by Theorem 6 .
In the induction step, assume that formula (3.29), as well as other conclusions in the Theorem 9 , holds true for $d$ dimensional case. In the $d+1$ dimension case, let $(\lambda, z)$ be one eigenpair of $A \in \mathbb{R}^{N^{(d+1)} \times N^{(d+1)}} . A$ is viewed as an $N \times N$ block matrix, and each block is of size $N^{d} \times N^{d}$. The eigenvector $z$ is partitioned to be

$$
z=\left(\begin{array}{c}
z_{1} \\
\vdots \\
z_{N}
\end{array}\right) \in \mathbb{R}^{N^{d+1}} \quad \text { with the } j \text {-th block of } z \text { being } z_{j} \in \mathbb{R}^{N^{d}}, \quad j=1, \ldots, N .
$$

Because $A z=\lambda z$, it gives the following difference equation

$$
\begin{equation*}
T z_{j-1}+(S-\lambda I) z_{j}+T z_{j+1}=0, \quad j=1,2, \ldots, N \tag{3.30}
\end{equation*}
$$

where $z_{0}=z_{N+1}=0, S$ and $T$ are $N^{d} \times N^{d}$ matrices defined in (3.28). $S$ has its main diagonal entries as $2(d+1)$. $S$ and $T$ can be diagonalized as $S=Q \Lambda_{S} Q^{T}$ and $T=Q \Lambda_{T} Q^{T}$ with the same orthogonal transformation matrix $Q$ and diagonal matrices $\Lambda_{S}$ and $\Lambda_{T}$.

We need to find out the eigenvalues of $S$ and $T$ in (3.30). Because $S$ is coincide to the $N^{d} \times N^{d}$ matrix $A$ defined in (3.28) with the only difference that the main diagonal entries of $S$ in (3.30) is $2(d+1)$, the eigenvalues for $S$ should be corrected as $\lambda_{k_{1}, k_{2}, \ldots, k_{d}}=2(d+1)-2 \sum_{i=1}^{d} \cos \left(\frac{k_{i} \pi}{N+1}\right)$, considering that all the conclusions in Theorem 9 are supposed to be true for $d$-D case. So the $k$-th diagonal block entries ( $k$ denotes the set $k_{1}, k_{2}, \ldots, k_{d}$ ) of $\Lambda_{S}$ and $\Lambda_{T}$ are

$$
\begin{aligned}
& \lambda_{S, k}=2(d+1)-2 \sum_{i=1}^{d} \cos \left(\frac{k_{i} \pi}{N+1}\right), \\
& \lambda_{T, k}=-1 .
\end{aligned}
$$

The $m$ th block component ( $m$ denotes the set $m_{1}, m_{2}, \ldots, m_{d}$ ) of column $k$ of $Q$ is

$$
q_{m}^{(k)}=\left(\frac{2}{N+1}\right)^{\frac{d}{2}} \prod_{i=1}^{d} \sin \left(\frac{m_{i} k_{i} \pi}{N+1}\right), \quad m_{i}, k_{i}=1,2, \ldots, N .
$$

Using (3.30), we have

$$
\begin{aligned}
& Q^{T} T z_{j-1}+Q^{T}(S-\lambda I) z_{j}+Q^{T} T z_{j+1}=0 \\
\Rightarrow \quad & \Lambda_{T} y_{j-1}+\left(\Lambda_{S}-\lambda I\right) y_{j}+\Lambda_{T} y_{j+1}=0, \quad \text { where } y_{j}=Q^{T} z_{j}, \quad j=1,2, \ldots, N .
\end{aligned}
$$

Both $\Lambda_{T}$ and $\left(\Lambda_{S}-\lambda I\right)$ are diagonal. So we have the following difference equation

$$
\lambda_{T, k} y_{k, j-1}+\lambda_{S, k} y_{k, j}+\lambda_{T, k} y_{k, j+1}=\lambda y_{k, j}, \quad j=1,2, \ldots, N .
$$

This problem can be rewritten in matrix form, when the index $k$ is fixed and all components of $y$ except the $k$ th one are 0 ,

$$
\left(\begin{array}{ccccc}
\lambda_{S, k} & \lambda_{T, k} & & & \\
\lambda_{T, k} & \lambda_{S, k} & \lambda_{T, k} & & \\
& \ddots & \ddots & \ddots & \\
& & \lambda_{T, k} & \lambda_{S, k} & \lambda_{T, k} \\
& & & \lambda_{T, k} & \lambda_{S, k}
\end{array}\right)\left(\begin{array}{c}
y_{k, 1} \\
y_{k, 2} \\
\vdots \\
y_{k, N-1} \\
y_{k, N}
\end{array}\right)=\lambda\left(\begin{array}{c}
y_{k, 1} \\
y_{k, 2} \\
\vdots \\
y_{k, N-1} \\
y_{k, N}
\end{array}\right) .
$$

By Theorem 7, the eigenpairs for such an eigenproblem are given as

$$
\begin{align*}
\lambda_{k, k_{d+1}} & =\lambda_{S, k}+2 \lambda_{T, k} \cos \left(\frac{k_{d+1} \pi}{N+1}\right) \\
& =2(d+1)-2 \sum_{i=1}^{d} \cos \left(\frac{k_{i} \pi}{N+1}\right)-2 \cos \left(\frac{k_{d+1} \pi}{N+1}\right)  \tag{3.31}\\
& =4 \sum_{i=1}^{d+1} \sin ^{2}\left(\frac{k_{i} \pi}{2(N+1)}\right) \\
y_{k, m_{d+1}}^{\left(k, k_{d+1}\right)} & =\sqrt{\frac{2}{N+1}} \sin \left(\frac{m_{d+1} k_{d+1} \pi}{N+1}\right), \quad m_{d+1}, k_{d+1}=1,2, \ldots, N
\end{align*}
$$

where $k$ denotes the set $\left\{k_{i}\right\}, m$ denotes $\left\{m_{i}\right\}, i=1,2, \ldots, d$. (3.31) shows that the term $2(d+1)$ is equal to the main diagonal entry of $A$.

Since $y_{j}=Q^{T} z_{j}, j=1,2, \ldots, N$, and because all components of $y$ except the $k$ th one are 0 , only the $k$ th entry of the $(d+1)$ th block of $y$ is nonzero. Hence, we have

$$
z_{m_{1}, m_{2}, \ldots, m_{d+1}}^{\left(k_{1}, k_{2}, \ldots, k_{d+1}\right)}=q_{m}^{k} y_{k, m_{d+1}}^{\left(k, k_{d+1}\right)}=\left(\frac{2}{N+1}\right)^{\frac{d+1}{2}} \prod_{i=1}^{d+1} \sin \left(\frac{m_{i} k_{i} \pi}{N+1}\right) .
$$

## 4 The Optimal Parameter for the Model Problem

Because many of the real world problems are large in scale, the matrix derived from discretizing the Poisson equation is large and sparse. Iterative methods provide efficient and economical ways to solve these linear large scale sparse systems. Jacobi, Gauss-Seidel, SOR, and SSOR are classic iterative methods, and the conjugate gradient method is an example of a modern iterative method. Preconditioning techniques are also available to improve the speed of convergence of these methods. Such techniques improve the spectral properties of these iteration matrix. As we have shown before, the spectral radius plays an important role in the convergence of the iterative method. Here, we explore more details about how to choose the optimal parameter $\omega$ in the SOR method. In fact, for the system matrix arising from the model problem, we will derive a formula for $\omega$ that minimizes the spectral radius of the SOR iteration matrix and is thus optimal.

In the model problem, the system matrix is $A=D-L-U$, where $D$ is the diagonal part of $A,-L$ the strictly lower triangular part, and $-U$ the strictly upper triangular part. Let $G_{J}$ and $G_{\omega}$ denote the iteration matrices of the Jacobi and SOR methods. Here, $G_{J}=I-D^{-1} A=D^{-1}(L+U)$, and $G_{\omega}=I-\left(\omega^{-1} D-L\right)^{-1} A=(D-\omega L)^{-1}[(1-\omega) D+\omega U]$.

Lemma 1. For the d-dimensional model problem ( $d \geq 1$ ), its system matrix $A=D-L-U$ defined in (3.28) satisfies (2.3), namely

$$
\begin{equation*}
\operatorname{det}\left(k D-\gamma L-\gamma^{-1} U\right)=\operatorname{det}(k D-L-U) \quad \text { for all } k, \gamma \in \mathbb{R} \backslash\{0\} \text {, } \tag{4.1}
\end{equation*}
$$

where $A$ is a block matrix, $D$ is diagonal of $A$, and $-L$ and $-U$ are the strictly lower and strictly upper triangular parts of $A$, respectively.

Proof. (3.28) shows how the system matrix in $d$-D case is constructed. After applying $k D-\gamma L-$ $\gamma^{-1} U$ to system matrix $A=D-L-U$, the matrix $S_{d-1}$ becomes a new matrix, denoted as $M_{d-1}$ such that

$$
M_{d-1}=\left[\begin{array}{ccccc}
2 k d & -\gamma^{-1} & & & \\
-\gamma & 2 k d & -\gamma^{-1} & & \\
& \ddots & \ddots & \ddots & \\
& & -\gamma & 2 k d & -\gamma^{-1} \\
& & & -\gamma & 2 k d
\end{array}\right]
$$

Each $S_{i}$ in (3.28) becomes

$$
M_{i}=\left[\begin{array}{ccccc}
M_{i+1} & \gamma^{-1} T & & & \\
\gamma T & M_{i+1} & \gamma^{-1} T & & \\
& \ddots & \ddots & \ddots & \\
& & \gamma T & M_{i+1} & \gamma^{-1} T \\
& & & \gamma T & M_{i+1}
\end{array}\right], \quad 1 \leq i \leq(d-2),
$$

where $T=-I$ are negative identity matrices with dimensions adjusted implicitly according to the size of $M_{i+1}$. We introduce notation $D_{i}$ to stand for the diagonal matrix of $S_{i},-L_{i}$ and $-U_{i}$ are the strictly lower and strictly upper triangular parts of $S_{i}$. Hence, $M_{i}=k D_{i}-\gamma L_{i}-\gamma^{-1} U_{i}$. Construct an invertible matrix $Q_{d-1}$ of size $N \times N$ that

$$
Q_{d-1}=\left[\begin{array}{lllll}
1 & & & & \\
& \gamma & & & \\
& & \gamma^{2} & & \\
& & & \ddots & \\
& & & & \gamma^{N-1}
\end{array}\right]
$$

Basing on that, a group of invertible block matrix $Q_{i}$ can be obtained such that

$$
Q_{i}=\left[\begin{array}{cccc}
Q_{i+1} & & & \\
& \gamma Q_{i+1} & & \\
& & \ddots & \\
& & & \gamma^{N-1} Q_{i+1}
\end{array}\right] \quad 1 \leq i \leq(d-2)
$$

It can be checked that the similarity transformation holds

$$
Q_{d-1}^{-1} M_{d-1} Q_{d-1}=Q_{d-1}^{-1}\left(k D_{d-1}-\gamma L_{d-1}-\gamma^{-1} U_{d-1}\right) Q_{d-1}=k D_{d-1}-L_{d-1}-U_{d-1} .
$$

Then consider the matrix $S_{d-1}$. Construct another invertible block matrix $Q_{d-2}$ that

$$
Q_{d-2}=\left[\begin{array}{llll}
Q_{d-1} & & & \\
& \gamma Q_{d-1} & & \\
& & \ddots & \\
& & & \gamma^{N-1} Q_{d-1}
\end{array}\right]
$$

Repeating the similarity transformation again, we have

$$
\begin{aligned}
& Q_{d-2}^{-1} M_{d-2} Q_{d-2}=Q_{d-2}^{-1}\left(k D_{d-2}-\gamma L_{d-2}-\gamma^{-1} U_{d-2}\right) Q_{d-2} \\
& =Q_{d-2}^{-1}\left[\begin{array}{ccccc}
M_{d-1} & \gamma^{-1} T & & & \\
\gamma T & M_{d-1} & \gamma^{-1} T & & \\
& \ddots & \ddots & \ddots & \\
& & & \gamma T & M_{d-1}
\end{array}\right] Q_{d-2} \\
& =\left[\begin{array}{ccccc}
Q_{d-1}^{-1} M_{d-1} & \gamma^{-1} Q_{d-1}^{-1} T & & \\
Q_{d-1}^{-1} T & \gamma^{-1} Q_{d-1}^{-1} M_{d-1} & \gamma^{-2} Q_{d-1}^{-1} T & \\
& \ddots & \ddots & \ddots \\
& & & \gamma^{-(N-2)} Q_{d-1}^{-1} T & \gamma^{-(N-1)} Q_{d-1}^{-1} M_{d-1}
\end{array}\right] Q_{d-2} \\
& =\left[\begin{array}{ccccc}
Q_{d-1}^{-1} M_{d-1} Q_{d-1} & T & & \\
T & & Q_{d-1}^{-1} M_{d-1} Q_{d-1} & T & \\
& & \ddots & \ddots & \ddots \\
& & & T & Q_{d-1}^{-1} M_{d-1} Q_{d-1}
\end{array}\right] \\
& =k D_{d-2}-L_{d-2}-U_{d-2} .
\end{aligned}
$$

After repeating $d-1$ times, we have an invertible matrix $Q$ such that

$$
Q^{-1}\left(k D-\gamma L-\gamma^{-1} U\right) Q=k D-L-U
$$

where $Q$ is a block diagonal matrix with its $i$ th diagonal block being $\gamma^{i-1} Q_{1},(1 \leq i \leq N)$. Therefore, $\operatorname{det}\left(k D-\gamma L-\gamma^{-1} U\right)=\operatorname{det}(k D-L-U)$, which is independent of $\gamma$ for all $\gamma \neq 0$ and for all $k$.

The case $d=1$ is a special case, however, in the 1-D case, the system matrix $A$ defined in (3.7) also satisfies (2.3). Having this, Theorem 3 and Theorem 4 can be applied to the system matrix $A$ of $d-\mathrm{D}$ model problem.

Theorem 10. Let $A=D-L-U$ be a matrix defined in (3.28). $G_{J}=I-D^{-1} A$ is the iteration matrix for Jacobi method. The eigenvalues of $G_{J}$ are

$$
\lambda_{k_{1}, k_{2}, \ldots, k_{d}}=1-\frac{2}{d} \sum_{i=1}^{d} \sin ^{2}\left(\frac{k_{i} \pi}{2(N+1)}\right), \quad k_{i}=1,2, \ldots, N .
$$

Here, $d \geq 1$ stands for the dimension of the problem.
Proof. For the $d$-D model problem, the eigenvalues of the system matrix are

$$
\lambda_{k_{1}, k_{2}, \ldots, k_{d}}^{\prime \prime}=4 \sum_{i=1}^{d} \sin ^{2}\left(\frac{k_{i} \pi}{2(N+1)}\right), \quad k_{i}=1,2, \ldots, N
$$

by Theorem 9 . Because the diagonal entries of $D$ for the $d$-D model problem are $2 d$, the eigenvalues of $D^{-1} A$ can be determined as

$$
\lambda_{k_{1}, k_{2}, \ldots, k_{d}}^{\prime}=\frac{2}{d} \sum_{i=1}^{d} \sin ^{2}\left(\frac{k_{i} \pi}{2(N+1)}\right), \quad k_{i}=1,2, \ldots, N
$$

Therefore, eigenvalues for $G_{J}=I-D^{-1} A$ are

$$
\lambda_{k_{1}, k_{2}, \ldots, k_{d}}=1-\frac{2}{d} \sum_{i=1}^{d} \sin ^{2}\left(\frac{k_{i} \pi}{2(N+1)}\right), \quad k_{i}=1,2, \ldots, N .
$$

Corollary 2. The spectral radius of $G_{J}$ is the same for any d-D model problem, namely

$$
\beta=\rho\left(G_{J}\right)=\cos \left(\frac{\pi}{N+1}\right) .
$$

Proof. According to Theorem 10, the spectral radius of $G_{J}$ is

$$
\beta=\rho\left(G_{J}\right)=1-\frac{2}{d} \sum_{i=1}^{d} \sin ^{2}\left(\frac{\pi}{2(N+1)}\right)=1-2 \sin ^{2}\left(\frac{\pi}{2(N+1)}\right)=\cos \left(\frac{\pi}{N+1}\right) .
$$

Theorem 11. Consider the model problem (1.1)-(1.2) with domain $\Omega \subset \mathbb{R}^{d}$ in $d \geq 1$ dimensions discretized by the finite difference method on a mesh with $N+2$ points in each coordinate direction and uniform mesh spacing $h=1 /(N+1)$. The optimal relaxation parameter for the $S O R$ method is then given by

$$
\begin{equation*}
\omega_{o p t}=\frac{2}{1+\sin \pi h}, \tag{4.2}
\end{equation*}
$$

which is independent of dimension $d$.


Figure 2: Comparison of different $\omega$ 's effect on iteration numbers.

Proof. Corollary 2, Theorem 9, formula (2.6) give

$$
\omega_{\text {opt }}=\frac{2}{1+\sqrt{1-\cos ^{2}\left(\frac{\pi}{N+1}\right)}}=\frac{2}{1+\sin \left(\frac{\pi}{N+1}\right)}=\frac{2}{1+\sin \pi h} .
$$

To confirm the optimality of this value of the relaxation parameter, we consider the 2-D model problem (3.14) with boundary condition (3.15) using a $65 \times 65$ mesh to partition the unit square $(0,1) \times(0,1)$. The function $f(x, y)$ is chosen as $f(x, y)=-2 \pi^{2} \cos (2 \pi x) \sin ^{2}(\pi y)-2 \pi^{2} \sin ^{2}(\pi x) \cos (2 \pi y)$. We solve this problem with the SOR method, with error tolerance $10^{-7}$ and maximum number of iterations allowed 1000 , and let the parameter $\omega$ vary from 1.5 to 2 . In this way, we can observe how the choice of $\omega$ affects the iteration. The optimal $\omega$ is also calculated according to (4.2), that is $\omega_{\text {opt }}=1.9078$, and its iteration number (206 times) is also obtained to compare with other $\omega$. In Figure 2, the smallest numbers of iteration is obtained when $\omega$ is near 1.9078.

## 5 Conclusions

There does not exist an explicit formula to compute the optimal relaxation parameter for SOR method in terms of properties of the system matrix $A$ of a general linear system. But the special structure of the system matrix $A$ resulting from the finite difference discretization of the classical model problem (1.1)-(1.2) allows for the derivation of such an explicit formula, based on the explicit determination of the spectrum of $A$. This result is well-known in two dimensions. Here, we presented a complete proof to extend this result to any dimension.

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[^0]:    *Department of Mathematics and Statistics, University of Maryland, Baltimore County, 1000 Hilltop Circle, Baltimore, MD 21250, \{shiming1,gobbert\}@math.umbc.edu

