CONJUGATE GRADIENT METHODS FOR TOEPLITZ SYSTEMS

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Abstract. In this expository paper, we survey some of the latest developments in using preconditioned conjugate gradient methods for solving Toeplitz systems. One of the main results is that the complexity of solving a large class of $n$-by-$n$ Toeplitz systems is reduced to $O(n \log n)$ operations as compared to $O(n \log^2 n)$ operations required by fast direct Toeplitz solvers. Different preconditioners proposed for Toeplitz systems are reviewed. Applications to Toeplitz-related systems arising from partial differential equations, queueing networks, signal and image processing, integral equations, and time series analysis are given.

Key words. Toeplitz matrices, preconditioners, preconditioned conjugate gradient methods, differential equations, signal and image processing, time series, queueing problems, integral equations

AMS subject classifications. 45E10, 62M10, 65F10, 65N22, 65P05, 68U10, 93E11

1. Introduction.

1.1. Background. An $n$-by-$n$ matrix $A_n$ is said to be Toeplitz if

\[
A_n = \begin{bmatrix}
a_0 & a_{-1} & \cdots & a_{2-n} & a_{1-n} \\
a_1 & a_0 & \cdots & a_{1-n} & a_{2-n} \\
\vdots & a_1 & a_0 & \cdots & \vdots \\
a_{n-2} & \cdots & a_1 & a_0 \\
a_{n-1} & a_{n-2} & \cdots & a_1 \end{bmatrix}
\]

i.e., $A_n$ is constant along its diagonals. The name Toeplitz originates from the work of Otto Toeplitz [180] in the early 1900s on bilinear forms related to Laurent series; see Grenander and Szegö [91] for details. We are interested in solving the Toeplitz system $A_n x = b$.

Toeplitz systems arise in a variety of applications in mathematics and engineering. In signal processing, solutions of Toeplitz systems are required in order to obtain the filter coefficients in the design of recursive digital filters; see Chui and A. Chan [61] and Haykin [99, pp. 167–168]. Time series analysis involves solutions of Toeplitz systems for the unknown parameters of stationary autoregressive models; see King et al. [125, pp. 368–379]. By using discrete time and spatial sampling of the domain and the sinc function $\sin(\pi x)/(\pi x)$ as a basis function for approximating the initial data, the numerical solutions of inverse heat problems can be obtained by solving a Toeplitz system; see Gilliam, Martin, and Lund [83]. Other applications involve solutions of partial differential equations, solutions of convolution-type integral equations, Padé approximations, and minimum realization problems in control theory; see Bunch [23] and the references therein. These applications have motivated both mathematicians and engineers to develop specific algorithms catering to solving Toeplitz systems. We will call these algorithms Toeplitz solvers.

Most of the early works on Toeplitz solvers were focused on direct methods. A straightforward application of the Gaussian elimination method will result in an algorithm of $O(n^3)$ complexity. However, since $n$-by-$n$ Toeplitz matrices are determined by only $(2n-1)$ entries rather than $n^2$ entries, it is expected that the solution of Toeplitz systems can be obtained in less than $O(n^3)$ operations. There are a number of Toeplitz solvers that decrease the complexity to
$O(n^2)$ operations; see, for instance, Levinson (1946) [134], Baxter (1961) [8], Trench (1964) [182], and Zohar (1974) [199]. These algorithms require the invertibility of the $(n-1)$-by-(n-1) principal submatrix of $A_n$. Around 1980, fast direct Toeplitz solvers of complexity $O(n \log^2 n)$ were developed; see, for instance, Brent, Gustavson, and Yun (1980) [20], Bitmead and Anderson (1980) [16], Morf (1980) [143], de Hoog (1987) [108], and Ammar and Gragg (1988) [4]. These algorithms require the invertibility of the $[n/2]$-by-$[n/2]$ principal submatrix of $A_n$.

The stability properties of these direct methods for symmetric positive-definite Toeplitz systems are discussed in Bunch [23]. It is noted that if $A_n$ has a singular or ill-conditioned principal submatrix, then a breakdown or near-breakdown can occur in these algorithms. Such breakdowns will cause numerical instabilities in subsequent steps of the algorithms and result in inaccurately computed solutions. The question of how to avoid breakdowns or near-breakdowns by skipping over singular submatrices or ill-conditioned submatrices has been studied extensively, and various algorithms have been proposed; see [58, 70, 89, 93, 102, 157, 175, 198]. In particular, T. Chan and Hansen (1992) [58] were the first to derive a look-ahead variant of the Levinson algorithm. The basic idea is to relax the inverse triangular decomposition slightly and to compute an inverse block factorization of the Toeplitz matrices with a block diagonal matrix instead of a scalar diagonal matrix. Other look-ahead extensions of fast Toeplitz solvers can be found in [77, 80, 95].

Recent research on using the preconditioned conjugate gradient method as an iterative method for solving Toeplitz systems has garnered much attention. One of the main important results of this methodology is that the complexity of solving a large class of Toeplitz systems can be reduced to $O(n \log n)$ operations as compared to the $O(n \log^2 n)$ operations required by fast direct Toeplitz solvers, provided that a suitable preconditioner is chosen under certain conditions on the Toeplitz operator. Besides the reduction of the arithmetic complexity, there are large classes of important Toeplitz matrices where the fast direct Toeplitz solvers are notoriously unstable, e.g., indefinite and certain non-Hermitian Toeplitz matrices. Therefore, iterative methods provide alternatives to solving these Toeplitz systems. In this paper, we will survey results for these iterative Toeplitz solvers and give some insight in how to construct effective preconditioners for them. Applications of these Toeplitz solvers to some practical problems will also be given.

1.2. Toeplitz matrices and circulant matrices. Let us begin by introducing the notation that will be used throughout the paper. Let $C_{2\pi}$ be the set of all $2\pi$-periodic continuous real-valued functions defined on $[-\pi, \pi]$. For all $f \in C_{2\pi}$, let

$$a_k = \frac{1}{2\pi} \int_{-\pi}^{\pi} f(\theta) e^{-ik\theta} d\theta, \quad k = 0, \pm 1, \pm 2, \ldots$$

be the Fourier coefficients of $f$. For all $n \geq 1$, let $A_n$ be the $n$-by-$n$ Toeplitz matrix with entries $a_{j,k} = a_{j-k}$, $0 \leq j, k < n$. The function $f$ is called the generating function of the sequence of Toeplitz matrices $A_n$; see Grenander and Szegö [91]. Since $f$ is a real-valued function, we have

$$a_{-k} = \bar{a}_k, \quad k = 0, \pm 1, \pm 2, \ldots$$

It follows that $A_n$ are Hermitian matrices. Note that when $f$ is an even function, the matrices $A_n$ are real symmetric. We emphasize that in practical applications, the functions $f$ are readily available. Typical examples of generating functions are the kernels of Wiener–Hopf equations (see Gohberg and Fel'dman [85, p. 82]), the functions which give the amplitude characteristics of recursive digital filters (see Chui and A. Chan [61]), the spectral density
functions in stationary stochastic processes (see Grenander and Szegő [91, p. 171]), and the point-spread functions in image deblurring (see Jain [116, p. 269]).

We will solve the systems \( A_n x = b \) by conjugate gradient methods. The convergence rate of the methods depends partly on how clustered the spectra of the sequence of matrices \( A_n \) are; see Axelsson and Barker [7, p. 24]. The clustering of the spectra of a sequence of matrices is defined as follows.

**Definition 1.1.** A sequence of matrices \( \{ A_n \}_{n=1}^{\infty} \) is said to have clustered spectra around 1 if for any given \( \epsilon > 0 \) there exist positive integers \( n_1 \) and \( n_2 \) such that for all \( n > n_1 \), at most \( n_2 \) eigenvalues of the matrix \( A_n - I_n \) have absolute value larger than \( \epsilon \).

For Toeplitz matrices, we note that there is a close relationship between the spectrum of \( A_n \) and its generating function \( f \).

**Theorem 1.1** (see Grenander and Szegő [91, pp. 63–65]). Let \( f \in \mathbb{C}_{2\pi} \). Then the spectrum \( \lambda(A_n) \) of \( A_n \) satisfies

\[
\lambda(A_n) \subseteq [f_{\min}, f_{\max}] \quad \forall n \geq 1,
\]

where \( f_{\min} \) and \( f_{\max} \) are the minimum and maximum values of \( f \), respectively. Moreover, the eigenvalues \( \lambda_j(A_n) \), \( j = 0, 1, \ldots, n - 1 \), are equally distributed as \( f(2\pi j/n) \); i.e.,

\[
\lim_{n \to \infty} \frac{1}{n} \sum_{j=0}^{n-1} \left[ g(\lambda_j(A_n)) - g \left( f \left( \frac{2\pi j}{n} \right) \right) \right] = 0
\]

for any continuous function \( g \) defined on \([-\pi, \pi]\).

The equal distribution of eigenvalues of Toeplitz matrices indicates that the eigenvalues will not be clustered in general. To illustrate this, consider the one-dimensional (1-d) discrete Laplacian matrix

\[
A_n = \text{tridiag}[-1, 2, -1].
\]

Its eigenvalues are given by

\[
\lambda_j(A_n) = 4 \sin^2 \left( \frac{\pi j}{n+1} \right), \quad 1 \leq j \leq n.
\]

For \( n = 32 \), the eigenvalues of \( A_n \) are depicted in Figure 1.1.

An \( n \)-by-\( n \) matrix \( C_n \) is said to be circulant if

\[
C_n = \begin{bmatrix}
c_0 & c_{-1} & \cdots & c_{2-n} & c_{1-n} \\
c_1 & c_0 & \cdots & c_{2-n} & c_{1-n} \\
\vdots & c_1 & c_0 & \cdots & \vdots \\
c_{n-2} & \cdots & \cdots & c_1 & c_{1-n} \\
c_{n-1} & c_{n-2} & \cdots & c_1 & c_0
\end{bmatrix},
\]

where \( c_{-k} = c_{n-k} \) for \( 1 \leq k \leq n - 1 \). Circulant matrices are diagonalized by the Fourier matrix \( F_n \); i.e.,

\[
C_n = F_n^* A_n F_n,
\]
where the entries of $F_n$ are given by

$$[F_n]_{j,k} = \frac{1}{\sqrt{n}} e^{2\pi i j k / n}, \quad 0 \leq j, k \leq n - 1,$$

and $\Lambda_n$ is a diagonal matrix holding the eigenvalues of $C_n$; see, for instance, Davis [69, p. 73]. We note that $\Lambda_n$ can be obtained in $O(n \log n)$ operations by taking the fast Fourier transform (FFT) of the first column of $C_n$. For the FFT algorithm, we refer to Cooley and Tukey [66]. In fact, the diagonal entries $\lambda_k$ of $\Lambda_n$ are given by

$$\lambda_k = \sum_{j=0}^{n-1} c_j e^{2\pi i j k / n}, \quad k = 0, \ldots, n - 1.$$  \hfill (1.5)

Once $\Lambda_n$ is obtained, the products $C_n y$ and $C_n^{-1} y$ for any vector $y$ can be computed by FFTs in $O(n \log n)$ operations using (1.4).

1.3. The conjugate gradient method for Toeplitz matrices. The conjugate gradient method is an iterative method for solving Hermitian positive-definite matrix systems. The algorithm of the method can be found in Golub and Van Loan [88, pp. 516–527]. In each iteration, it requires two inner products of $n$-vectors and one multiplication of the coefficient matrix with an $n$-vector. Storage for four temporary $n$-vectors is needed but there is no explicit storage required for the coefficient matrix.

Let $f \in C_{2\pi}$. For simplicity, let us assume for the moment that $f$ is positive; i.e., $f_{\min} > 0$. Then by (1.2) $A_n$ are positive definite for all $n$. Consider applying the conjugate gradient method to solve these symmetric positive-definite Toeplitz systems $A_n x = b$. In each iteration, besides the two inner products required, one matrix-vector multiplication $A_n y$ is also needed. That can be computed by FFTs by first embedding $A_n$ into a $2n$-by-$2n$ circulant matrix, i.e.,

$$\begin{pmatrix} A_n & \times \\
\times & A_n \end{pmatrix} \begin{bmatrix} y \\
0 \end{bmatrix} = \begin{bmatrix} A_n y \\
0 \end{bmatrix}$$

(see Strang [172]), and then carrying out the multiplication by using the decomposition as in (1.4). The matrix-vector multiplication thus requires $O(2n \log(2n))$ operations. It follows that the total number of operations per iteration is of $O(n \log n)$ operations. As for the storage required, besides the four temporary $n$-vectors, we need an extra $2n$-vector for storing the eigenvalues of the embedded circulant matrix given in (1.6).

The convergence rate of the conjugate gradient method has been well studied; see Axelsson and Barker [7, p. 24]. It depends on the condition number of the matrix $A_n$ and how clustered the spectrum of $A_n$ is. If the spectrum is not clustered, as is usually the case for Toeplitz matrices (cf. Theorem 1.1), a good estimate of the convergence rate is given in terms of the largest and smallest eigenvalues of $A_n$. Using (1.2), this estimate can be expressed as

$$\frac{\|e_q\|_{A_n}}{\|e_0\|_{A_n}} < 2 \left( \frac{\sqrt{f_{\max}} - \sqrt{f_{\min}}}{\sqrt{f_{\max}} + \sqrt{f_{\min}}} \right)^q,$$

where $e_q$ is the error vector at the $q$th iteration and $\|x\|_{A_n}^2 = x^* A_n x$. This indicates that the rate of convergence is linear. Thus, the method will converge in a constant number of iterations, and hence the complexity of solving the Toeplitz system is $O(n \log n)$. However, we remark that if $f_{\max}/f_{\min}$ is large, the constant in the operation count will be large and hence the convergence will be very slow.
One way to speed up the convergence rate of the method is to precondition the Toeplitz system. Thus, instead of solving \( A_n x = b \), we solve the preconditioned system

\[
P_n^{-1} A_n x = P_n^{-1} b.
\]

The matrix \( P_n \), called the preconditioner, should be chosen according to the following criteria.

- \( P_n \) should be constructed within \( O(n \log n) \) operations.
- \( P_n v = y \) should be solved in \( O(n \log n) \) operations.
- The spectrum of \( P_n^{-1} A_n \) should be clustered.

The first two criteria are to keep the operation count per iteration within \( O(n \log n) \), as that is the count for the nonpreconditioned system. The third criterion comes from the fact that the more clustered the eigenvalues are, the faster the convergence of the method will be; see, for instance, [139, pp. 249–251] and [7, pp. 27–28]. If \( P_n^{-1} A_n \) has a clustered spectrum as defined in Definition 1.1, then the conjugate gradient method, when applied to solving the preconditioned system (1.7), converges superlinearly for large \( n \); see [49]. More precisely, for any given \( \epsilon > 0 \), there exists a constant \( c(\epsilon) > 0 \) such that the error vector \( e_q \) of the preconditioned conjugate gradient method at the \( q \)th iteration satisfies

\[
\frac{||e_q||}{||e_0||} \leq c(\epsilon) \epsilon^q,
\]

where

\[
|||v|||^2 = v^* P_n^{-1/2} A_n P_n^{-1/2} v.
\]

The main aim of this paper is to review different preconditioners developed for Toeplitz systems that satisfy the three criteria mentioned earlier. We will also study applications to Toeplitz-related systems arising from partial differential equations, queueing networks, signal and image processing, integral equations, and time series analysis. For simplicity, we will drop the subscripts on matrices when their dimensions are apparent. The outline of the paper is as follows. In \( \S 2 \), we survey the use of circulant matrices as preconditioners for Toeplitz matrices. In \( \S 3 \), other useful and successful noncirculant preconditioners for Toeplitz matrices are also considered. Applications of preconditioned conjugate gradient methods for Toeplitz-related systems are discussed in \( \S 4 \). Finally, concluding remarks are given in \( \S 5 \).

2. Circulant preconditioners for Toeplitz systems.

2.1. Circulant preconditioners. In 1986, Strang [172] and Olkin [155] independently proposed the use of circulant matrices to precondition Toeplitz matrices in conjugate gradient iterations. Part of their motivation was to exploit the fast inversion of circulant matrices. Numerical results in [173, 155] suggest that the method converges very fast for a wide range of Toeplitz matrices. This was later proven theoretically in [49] and in other papers for other circulant preconditioners. In this subsection, we will give a brief account of these developments.

With circulant matrices as preconditioners, in each iteration, we must solve a circulant system. From (1.4), we see that circulant matrices can be diagonalized by discrete Fourier matrices, and hence the inversion of \( n \)-by-\( n \) circulant systems can be done in \( O(n \log n) \) operations by using FFTs of size \( n \). In contrast, by (1.6), we see that the cost of computing \( Ay \), which is also required in each iteration whether the system is preconditioned or not, is done by using FFTs of size \( 2n \). Notice that if FFT is used to compute the discrete Fourier transform of a \( 2n \)-vector for which the even discrete Fourier transform components are already known, then the cost is the same as carrying out a length \( n \) FFT; see Linzer [136] for details. Thus,
the cost per iteration of the circulant preconditioned conjugate gradient method is roughly 1.25 times that required by the method without using preconditioners. We remark that Huckle [113] has also discussed different ways to reduce the number of FFTs in the iterative scheme, even when \( n \) is not a power of 2. In particular, he proposed a way to compute \( Ay \) such that the computational cost per iteration of the preconditioned system is nearly the same as that required by the nonpreconditioned system.

We emphasize that the use of circulant matrices as preconditioners for Toeplitz systems allows the use of FFT throughout the computations, and FFT is highly parallelizable and has been efficiently implemented on multiprocessors [2, p. 238] and [174]. Since conjugate gradient methods are also easily parallelizable [12, p. 165], the circulant preconditioned conjugate gradient method is well adapted for parallel computing.

We remark that circulant approximations to Toeplitz matrices have been considered and used for some time in signal processing (e.g., [158] and [171, pp. 75–86]), time series analysis (e.g., [21, p. 133] and [169]), and image processing (e.g., [13], [123, p. 147], and [5, p. 136]). However, in these applications, the circulant approximations thus obtained were used to replace the given Toeplitz matrices in subsequent computations. In contrast, circulant approximations are used here only as preconditioners for Toeplitz systems and the solutions to the Toeplitz systems are unchanged. In the following, we review some successful circulant preconditioners proposed for Toeplitz matrices.

2.1.1. Strang’s preconditioner. The first circulant preconditioner was proposed by Strang [172] in 1986 and is defined to be the matrix that copies the central diagonals of \( A \) and reflects them around to complete the circulant requirement. For \( A \) given by (1.1), the diagonals \( s_j \) of the Strang preconditioner \( S = [s_{k-l}]_{0 \leq k, l < n} \) are given by

\[
\begin{align*}
    s_j &= \begin{cases} 
    a_j, & 0 < j \leq \lfloor n/2 \rfloor, \\
    a_{j-n}, & \lfloor n/2 \rfloor < j < n, \\
    s_{n-j}, & 0 < -j < n.
    \end{cases}
\end{align*}
\]

(2.1)

One of the interesting properties of \( S \) is that it minimizes

\[
||C - A||_1 \quad \text{and} \quad ||C - A||_\infty
\]

over all Hermitian circulant matrices \( C \); see [27]. The spectra of these circulant preconditioned matrices have been analyzed by R. Chan and Strang [49].

THEOREM 2.1 (see R. Chan and Strang (1989) [49]). Let \( f \) be an even positive function in the Wiener class; i.e., its Fourier coefficients are absolutely summable,

\[
\sum_{k=0}^{\infty} |a_k| < \infty.
\]

Let \( A \) be generated by \( f \). Then the spectra of \( S^{-1}A \) are clustered around 1 for large \( n \).

The main idea of their proof is to use an orthogonal transformation to transform \( S - A \) into a Hankel matrix. Then Nehari’s theorem [62, p. 120] is used to show that the limiting Hankel operator is compact. Using the theory of collectively compact sets of operators [6, pp. 65–70], the spectra of the finite Hankel matrices are then shown to be clustered. However, this proof cannot be readily generalized to real-valued \( f \), i.e., to Hermitian Toeplitz systems. Thus R. Chan in [27] developed a purely linear algebra technique to extend the results in Theorem 2.1. The approach is to decompose, for a given \( \epsilon > 0 \), the matrix \( S - A \) into a sum of two matrices \( L \) and \( V \), where rank \( L \leq c(\epsilon) \) and \( ||V||_2 \leq \epsilon \), and then apply Cauchy’s interlace theorem [88].
THEOREM 2.2 (see R. Chan (1989) [27]). Let \( f \) be a positive function in the Wiener class. Then the spectra of \( S^{-1}A \) are clustered around 1 for large \( n \).

Using standard error analysis of the conjugate gradient method, we can then show that the convergence rate of the method is superlinear; see (1.8). If extra smoothness conditions are imposed on the generating function \( f \), we can get more precise estimates on how \( \| \|e_{q+1}\| / \|e_q\| \) in (1.8) goes to zero.

THEOREM 2.3 (see Trefethen (1990) [181], Ku and Kuo (1993) [130, 131]). Suppose \( f \) is a rational function of the form \( f = p/q \), where \( p \) and \( q \) are polynomials of degrees \( \mu \) and \( \nu \), respectively. Then the number of outlying eigenvalues of \( S^{-1}A \) is exactly equal to \( 2 \max\{\mu, \nu\} \). Hence, the method converges in at most \( 1 + 2 \max\{\mu, \nu\} \) steps for large \( n \). If, however, \( f(z) = \sum_{j=0}^{\infty} a_j z^j \) is analytic only in a neighborhood of \( |z| = 1 \), then there exist constants \( c \) and \( 0 \leq r < 1 \) such that

\[
\frac{\|e_{q+1}\|}{\|e_q\|} \leq c^r q^{\lambda/4} q^{\lambda/2}.
\]

(2.2)

The idea of Trefethen's proof is to use rational approximation to bound the singular values of the Hankel matrix considered in the proof of Theorem 2.1. It follows from (2.2) that

\[
\frac{\|e_{q+1}\|}{\|e_q\|} \approx cr^q \rightarrow 0.
\]

For generating functions \( f \) with Fourier coefficients \( a_j \) decaying at a slower rate, we have the following two theorems.

THEOREM 2.4 (see R. Chan (1989) [27]). Let \( f \) be a \( \nu \)-times differentiable function with \( f^{(\nu)} \in L^1[-\pi, \pi] \), where \( \nu > 1 \) (i.e., \( |a_j| \leq \bar{c}/j^{\nu+1} \) for some constant \( \bar{c} \)). Then there exists a constant \( c \) that depends only on \( f \) and \( \nu \), such that for large \( n \),

\[
\frac{\|e_{2q}\|}{\|e_0\|} \leq \frac{c^q}{((q - 1)!)^{2\nu - 2}}.
\]

Theorem 2.4 was proven by using Cauchy's interlace theorem [88]. R. Chan and Yeung later used Jackson's theorem [60] in approximation theory to prove a stronger result than that in Theorem 2.4.

THEOREM 2.5 (see R. Chan and Yeung (1992) [55]). Let \( f \) be a Lipschitz function of order \( \nu \), \( 0 < \nu \leq 1 \), or \( f \) has a continuous \( \nu \)th derivative, \( \nu \geq 1 \). Then there exists a constant \( c \) that depends only on \( f \) and \( \nu \), such that for large \( n \),

\[
\frac{\|e_{2q}\|}{\|e_0\|} \leq \prod_{k=2}^{\nu} \frac{c \log^2 k}{k^{2\nu}}.
\]

Theorems 2.1–2.5 give the rate at which the error goes to zero in terms of the rate of decay of \( |a_j| \). To see if solving Toeplitz systems by preconditioned conjugate gradient methods is more efficient than using fast direct Toeplitz solvers, Linzer [137] performed tests for Toeplitz matrices with different condition numbers and coefficients with different decaying rates. His results show that the iterative methods have the edge if \( |a_j| \) decays like \( O(j^{-0.5}) \) or faster for matrices with condition number about 10, and the rate of decay required increases to \( O(j^{-2}) \) when the condition number is about 10^5.

We also remark that Huckle [114] has recently compared the number of floating point operations for iterative methods with those of direct Toeplitz solvers and superfast Toeplitz solvers. He derived an upper bound of the number of PCG iterations such that the iterative
method will be better than direct and superfast Toeplitz solvers. His finding shows that for positive-definite Toeplitz systems, PCG methods are competitive for large matrices with a small number of PCG iterations. In the indefinite or near-singular case, iterative methods may give a higher accuracy. In the unsymmetric case, only classical \( O(n^2) \) direct Toeplitz solvers are available, and therefore the iterative methods will have the edge if a good preconditioner can be found; see §2.3.

### 2.1.2. T. Chan's Preconditioner

For an \( n \)-by-\( n \) Toeplitz matrix \( A \), T. Chan's circulant preconditioner \( c(A) \) is defined to be the minimizer of

\[
\| C - A \|_F
\]

over all \( n \)-by-\( n \) circulant matrices \( C \); see T. Chan (1988) [57]. Here \( \| \cdot \|_F \) denotes the Frobenius norm. In [57], the matrix \( c(A) \) is called an optimal circulant preconditioner because it minimizes (2.3). The \( j \)-th diagonals of \( c(A) \) are shown to be equal to

\[
c_j = \begin{cases} 
(n - j)a_j + j a_{j-n} \\
\frac{c_{n+j}}{n}
\end{cases}, \quad 0 \leq j < n, \\
0 < -j < n,
\]

which are just the average of the diagonals of \( A \), with the diagonals being extended to length \( n \) by a wrap-around. By using (1.5) and (2.4), we see that the eigenvalues \( \lambda_k(c(A)) \) of \( c(A) \) are given by

\[
\lambda_k(c(A)) = \sum_{j=-n+1}^{n-1} a_j \left( 1 - \frac{|j|}{n} \right) e^{2\pi i j k / n}, \quad k = 0, \ldots, n - 1.
\]

As for the performance of \( c(A) \) as a preconditioner for \( A \), R. Chan [28] proved that under the Wiener class assumptions in Theorem 2.2 (i.e., \( f \) is a positive function with absolutely summable Fourier coefficients), the spectra of \( c(A)^{-1}A \) and \( S^{-1}A \) are asymptotically the same as \( n \) tends to infinity; i.e., \( \lim_{n \to \infty} ||c(A)^{-1}A - S^{-1}A||_2 = 0 \). Hence, \( c(A) \) works as well for Wiener class functions as \( S \) does. Using Weierstrass's theorem to approximate \( 2\pi \)-periodic continuous generating functions by Wiener class functions, we have the following theorem.

**Theorem 2.6** (see R. Chan and Yeung (1992) [53]). *Let \( f \) be a positive function in \( C_{2\pi} \). Then the spectra of \( c(A)^{-1}A \) are clustered around \( 1 \) for large \( n \).*

However, the Weierstrass approach used in proving this theorem does not work for Strang's preconditioner. From Theorem 2.5, we see that the class of generating functions where Strang's preconditioner works is the class of \( 2\pi \)-periodic Lipschitz continuous functions. We will discuss this discrepancy in §2.2.

When \( A \) is not a Toeplitz matrix, the circulant minimizer \( c(A) \) of (2.3) can still be obtained easily by taking the arithmetic average of the entries of \( A \); i.e., its diagonals are given by

\[
c_\ell = \frac{1}{n} \sum_{j-k=\ell (\text{mod} n)} a_{j,k}, \quad \ell = 0, \ldots, n - 1;
\]

see [183]. Therefore, T. Chan's preconditioner is particularly useful in solving non-Toeplitz systems arising from the numerical solutions of elliptic partial differential equations [31] and Toeplitz least squares problems arising from signal and image processing [42, 43, 44, 59, 97, 152, 153, 159]. Convergence results for T. Chan's preconditioner have been established for these problems; see §4.

Another interesting spectral property of \( c(A) \) is that if \( A \) is positive definite, then \( c(A) \) is also positive definite. In fact, for an arbitrary Hermitian matrix \( A \), we have

\[
\lambda_{\text{min}}(A) \leq \lambda_{\text{min}}(c(A)) \leq \lambda_{\text{max}}(c(A)) \leq \lambda_{\text{max}}(A);
\]
see Tyrtyshnikov [183] and R. Chan, Jin, and Yeung [38]. We remark that Strang’s preconditioner does not satisfy (2.7) even for Toeplitz matrices; see R. Chan and Yeung [54]. In addition, R. Chan and Wong [51] recently proved that for some Toeplitz matrices $A$, T. Chan’s preconditioner $c(A)$ minimizes $\kappa(C^{-1}A)$ over all $n$-by-$n$ nonsingular circulant matrices $C$.

Huckle in 1992 [111] proposed a preconditioner that is an extension of T. Chan’s preconditioner. Let $1 \leq p \leq n$. Huckle’s circulant preconditioner $H$ is defined to be the circulant matrix with eigenvalues

$$\lambda_k(H) = \sum_{j=-p+1}^{p-1} a_j \left(1 - \frac{|j|}{p}\right) e^{2\pi i j k / n}, \quad k = 0, \ldots, n - 1;$$

cf. (2.5). Thus, when $p = n$, it is the T. Chan preconditioner. Besides Wiener class functions, Huckle [110] has shown that $H$ also works for generating functions with Fourier coefficients $a_k$ that satisfy

$$\sum_{k=-\infty}^{\infty} |k||a_k|^2 < \infty.$$

2.1.3. Preconditioners by embedding. Let the Toeplitz matrix $A$ be embedded into a $2n$-by-$2n$ circulant matrix

$$(2.9) \quad \begin{bmatrix} A & B^* \\ B & A \end{bmatrix}.$$ 

R. Chan’s circulant preconditioner $R$ is defined as $R = A + B$; see R. Chan (1989) [27]. We remark that T. Chan’s and Huckle’s preconditioners to $A$ are just equal to R. Chan’s preconditioner obtained from the Toeplitz matrix with diagonals $(1 - |j|/n)a_j$ and $\max(0, 1 - |j|/p)a_j$, respectively; i.e., the diagonals $a_j$ of $A$ are damped by the factor $(1 - |j|/n)$ and $\max(0, 1 - |j|/p)$, respectively; see (2.5) and (2.8).

Using the embedding (2.9), Ku and Kuo (1992) [128] constructed four different preconditioners $K_i$, $1 \leq i \leq 4$, based on different combinations of the matrices $A$ and $B$. They are

$$K_1 = A + B, \quad K_2 = A - B, \quad K_3 = A + JB, \quad K_4 = A - JB,$$

where $J$ is the $n$-by-$n$ anti-identity (reversal) matrix. We note that $K_2$, $K_3$, and $K_4$ are not circulant matrices.

2.1.4. Preconditioners by minimization of norms. Besides using the minimizer of $\|C - A\|_F$ as preconditioners for Toeplitz systems, minimizers of other approximations have also been proposed and used. For instance, Tyrtyshnikov’s circulant preconditioner $T$ (1992) [183] is defined to be the minimizer of

$$\|I - C^{-1}A\|_F$$

over all nonsingular circulant matrices $C$. In [183], $T$ is called the superoptimal circulant preconditioner because it minimizes (2.10) instead of (2.3) and is shown to be equal to

$$T = c(AA^*)c(A)^{-1};$$

see also [38]. Tyrtyshnikov [183] showed that for a general positive-definite Toeplitz matrix $A$, $T$ is also positive definite; cf. (2.7). Tismenetsky [179] and Sayed and Kailath [165] independently proposed the same preconditioner.
Table 2.1

<table>
<thead>
<tr>
<th>n</th>
<th>I</th>
<th>S</th>
<th>c(A)</th>
<th>R</th>
<th>K₂</th>
<th>H</th>
<th>T</th>
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<td>5</td>
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<td>8</td>
</tr>
</tbody>
</table>

Huckle [112] has considered the minimizer \( M \) by minimization of

\[ \| I - C^{-1/2} AC^{-1/2} \|_F \]

over all nonsingular circulant matrices \( C \). The constructions of \( T \) and \( M \) require \( O(n \log n) \) operations; see R. Chan, Jin, and Yeung [38], Tismenetsky [179], and Huckle [112], respectively.

Finally, we compare the performance of these preconditioners \( (R, \{K_i\}_{i=1}^4, T, \text{ and } M) \) with Strang’s and T. Chan’s preconditioners. It has been proven in R. Chan [27], R. Chan, Jin, and Yeung [39], Ku and Kuo [128], and Huckle [112] that under the same Wiener class assumptions, these circulant preconditioned systems have spectra that are asymptotically the same as Strang’s and T. Chan’s circulant preconditioned systems. In particular, all these preconditioned systems converge at the same rate for large \( n \).

In the following, we illustrate the effectiveness of circulant preconditioners for Toeplitz systems by a numerical example. We use the continuous generating function

\[ f(\theta) = \theta^4 + 1, \quad -\pi \leq \theta \leq \pi, \]

in the test. Table 2.1 shows the numbers of iterations required to solve nonpreconditioned systems \( Ax = b \) and circulant preconditioned systems \( C^{-1} Ax = C^{-1} b \) for different preconditioners. The right-hand side vector \( b \) is the vector of all ones. The zero vector is the initial guess. The stopping criterion is when the residual vector \( r_q \) at the \( q \)th iteration satisfies \( \|r_q\|_2/\|r_0\|_2 < 10^{-7} \). All computations were done by Matlab.

In the table, \( I \) denotes that no preconditioner is used and \( H \) is Huckle’s preconditioner with \( p = n/2 \) (see (2.8)). We see from Table 2.1 that the number of iterations required for convergence for nonpreconditioned systems is much greater than those for circulant preconditioned systems. Figure 2.1 depicts the spectra of the nonpreconditioned matrix and the circulant preconditioned matrices for different circulant preconditioners. We note that the spectra of the circulant preconditioned matrices are indeed clustered around 1.

2.2. Circulant preconditioners from kernel functions. A unifying approach for constructing circulant preconditioners is given in R. Chan and Yeung [54], where it is shown that most of the above-mentioned circulant preconditioners can be derived by using the convolution products of some well-known kernels with the generating function \( f \). For example, Strang’s and T. Chan’s circulant preconditioners are constructed by using the Dirichlet and Fejér kernels, respectively. To see this, we start by noting that the eigenvalues of Strang’s preconditioner \( S \) are given by

\[ \lambda_j(S) = (\hat{D}_{1/2} * f)\left(\frac{2\pi j}{n}\right), \quad 0 \leq j < n, \]
where the convolution of the Dirichlet kernel with $f$ is given by

\[(\hat{D}_{\frac{n}{2}} * f)(\theta) = \frac{1}{2\pi} \int_{-\pi}^{\pi} \hat{D}_{\frac{n}{2}}(\theta - \phi) f(\phi) d\phi\]

and

\[\hat{D}_k(\theta) = \frac{\sin(k \cdot \frac{1}{2} \theta)}{\sin \frac{1}{2} \theta}, \quad k = 1, 2, \ldots.
\]

The eigenvalues of T. Chan's preconditioner $c(A)$ are given by

\[\lambda_j(c(A)) = (\hat{F}_n * f) \left( \frac{2\pi j}{n} \right), \quad 0 \leq j < n,
\]

where the Fejér kernels are given by

\[\hat{F}_k(\theta) = \frac{1}{k} \left( \frac{\sin \frac{k}{2} \theta}{\sin \frac{1}{2} \theta} \right)^2, \quad k = 1, 2, \ldots.
\]

Similarly, the eigenvalues of R. Chan's preconditioner $R$ and Huckle's preconditioner $H$ are given by

\[\lambda_j(R) = (\hat{D}_{n-1} * f) \left( \frac{2\pi j}{n} \right), \quad 0 \leq j < n
\]
and

$$\lambda_j(H) = (\hat{F}_p * f) \left( \frac{2\pi j}{n} \right), \quad 0 \leq j < n,$$

respectively. The idea can be applied to design other circulant preconditioners $C$ from kernels $\hat{C}_n$ such as the von Hann kernel, Hamming kernel, and Bernstein kernel, which are commonly used in function theory [190], signal processing [96], and time series analysis [160]. In [54], several circulant preconditioners were constructed using this approach.

We remark that the convolution products of these kernels with $f$ are just smooth approximations of the generating function $f$ itself. R. Chan and Yeung proved that if the convolution product is a good approximation of $f$, then the correspondingly constructed circulant matrix will be a good preconditioner.

**THEOREM 2.7** (see R. Chan and Yeung (1992) [54]). Let $f \in C_{2\pi}$ be positive. Let $\hat{C}_n$ be a kernel such that $\hat{C}_n * f$ tends to $f$ uniformly on $[-\pi, \pi]$. If $C$ is the circulant matrix with eigenvalues given by

$$\lambda_j(C) = (\hat{C}_n * f) \left( \frac{2\pi j}{n} \right), \quad 0 \leq j < n,$$

then the spectra of $C^{-1}A$ are clustered around 1 for large $n$.

We mentioned in §2.1 that Strang's and T. Chan's preconditioners are fundamentally different, in that Strang's preconditioner works for Lipschitz continuous functions (Theorem 2.5) while T. Chan's preconditioner works for $2\pi$-periodic continuous functions (Theorem 2.6). This can be explained by the associations of Strang's preconditioner with the Dirichlet kernel and T. Chan's preconditioner with the Fejér kernel. It is well known in Fourier analysis that if $f$ is $2\pi$-periodic continuous (or, respectively, Lipschitz continuous), then the convolution product of the Fejér kernel (or, respectively, the Dirichlet kernel) with $f$ will converge to $f$ uniformly on $[-\pi, \pi]$; see Walker [190, pp. 79, 52].

In addition, it is interesting to note that for a piecewise continuous function $f$, the convolution product with the Fejér kernel will no longer converge to $f$ uniformly on $[-\pi, \pi]$. Therefore, for generating functions that are only piecewise continuous, we don't expect the spectra of $c(A)^{-1}A$ to be clustered around 1.

**THEOREM 2.8** (see Yeung and R. Chan (1993) [195]). Let $f$ be a nonnegative piecewise continuous function on $[-\pi, \pi]$. Then for any given $\epsilon > 0$, the number of eigenvalues of $c(A)^{-1}A$ that lie outside the interval $(1 - \epsilon, 1 + \epsilon)$ is at least of $O(\log n)$ for $n$ sufficiently large. If, moreover, $f$ is strictly positive, then the number of outlying eigenvalues is exactly of $O(\log n)$.

The theorem is established by noting that $A - c(A)$ is orthogonally similar to a Hankel matrix. Then Widom's theorem [192], which gives an estimate of the eigenvalues of the Hilbert matrices, is used to estimate the number of outlying eigenvalues of $A - c(A)$. Numerical examples are given in [195] to verify that the convergence rate of the method will no longer be superlinear in general. In fact, the numbers of iterations required for convergence do increase like $O(\log n)$. These results have recently been extended by Tyryshnikov [184, 186]. In [184], he established a generalized Szegő theorem that if $f$ is in $L_2$, then the singular values of $A_n$ are equally distributed (in a generalized sense) as $|f(x)|$ (cf. (1.3)). He then used the result to prove that if $f$ is in $L_2$, then the number of outlying eigenvalues of the preconditioned system grows no more than $o(n)$. In [186], he further extended the results of circulant preconditioners for products of Toeplitz matrices.

### 2.3. Non-Hermitian-type Toeplitz systems

In this subsection, we study Toeplitz matrices $A$ generated by complex-valued functions. We note that such $A$ are complex non-Hermitian...
matrices. In general the fast direct Toeplitz solvers are not applicable, and neither is the conjugate gradient method when applied to the system $Ax = b$. For such matrices $A$, one can apply the conjugate gradient method to the normal equations $A^*Ax = A^*b$. Another way of solving non-Hermitian Toeplitz systems is to employ some CG-like method [78] such as restarted GMRES [164] or TFQMR [76]. To speed up the CG-like methods, we can choose a matrix $C$ such that the singular values of the preconditioned matrices $C^{-1}A$ are clustered.

Let us begin with skew-Hermitian-type Toeplitz matrices, i.e., Toeplitz matrices of the form

$$
A = \begin{bmatrix}
  a_0 & -\bar{a}_1 & \cdots & -\bar{a}_{n-2} & -\bar{a}_{n-1} \\
  a_1 & a_0 & -\bar{a}_1 & \cdots & -\bar{a}_{n-2} \\
  \vdots & a_1 & a_0 & \ddots & \vdots \\
  a_{n-2} & \ddots & \ddots & -\bar{a}_1 \\
  a_{n-1} & a_{n-2} & \cdots & a_1 & a_0
\end{bmatrix},
$$

where $a_0$ is a real number. Obviously, $A = a_0I + A_S$, where $I$ is the identity matrix and $A_S$ is a skew-Hermitian Toeplitz matrix. These systems or low-rank perturbations of such systems often appear in solving discretized hyperbolic differential equations; see Buckley [22] and Holm gren and Otto [104, 105].

In [35], R. Chan and Jin used R. Chan's circulant preconditioner $R$ defined in §2.1.3 to precondition these skew-Hermitian Toeplitz matrices. Under the Wiener class assumptions on the entries of the first column of $A$, they proved that the singular values of $R^{-1}A$ are clustered around 1.

For Toeplitz matrices $A$ generated by a complex-valued function, R. Chan and Yeung [56] have proven that if the generating function is $2\pi$-periodic continuous with no zeros on $[-\pi, \pi]$, then the spectra of the iteration matrices $(c(A)^{-1}A)^*c(A)^{-1}A$ are clustered around 1. From this they showed that if the condition number $\kappa(A)$ of $A$ is of $O(n^\alpha)$, $\alpha > 0$, then the number of iterations required for convergence is at most $O(\alpha \log n)$. Hence the total complexity for solving non-Hermitian-type Toeplitz systems is of $O(n \log^2 n)$. When $\alpha = 0$, i.e., $A$ is well conditioned, the method converges in $O(1)$ steps and the complexity is reduced to $O(n \log n)$.

Numerical results in [56] show that the requirements on $f$, namely, that $f$ has no zeros and $\kappa(A) = O(n^\alpha)$, are indispensable in order to get the said convergence rate. Moreover, these two conditions are mutually exclusive. For instance, if $f(\theta) = e^{i\theta}$, then $f$ has no zeros on $[-\pi, \pi]$ but $A$ is singular for all $n$. On the other hand, if $f(\theta) = 4 \sin^2 \theta$, then $A$ is just the 1-d discrete Laplacian with $\kappa(A) = O(n^2)$.

We remark that Ku and Kuo [130, 131, 132] have also considered solving nonsymmetric Toeplitz matrix systems by the preconditioned conjugate gradient method. In their papers, $A$ is assumed to be generated by complex-valued rational functions in the Wiener class, which happens to be a subclass of the class of $2\pi$-periodic continuous functions considered in [56].

### 2.4. \(\omega\)-circulant preconditioners

Circulant matrices belong to the class of \(\omega\)-circulant matrices, which are defined as follows.

**Definition 2.1.** Let $\omega = e^{i\theta_0}$ with $\theta_0 \in [-\pi, \pi]$. An $n$-by-$n$ matrix $W$ is said to be an \(\omega\)-circulant matrix if it has the spectral decomposition

$$W = \Omega^*F^* \Lambda F \Omega.$$

Here $\Omega = \text{diag}[1, \omega^{-1/n}, \ldots, \omega^{-(n-1)/n}]$ and $\Lambda$ is a diagonal matrix containing the eigenvalues of $W$. 
Notice that \( \{ \omega \} \)-circulant matrices are Toeplitz matrices with the first entry of each row obtained by multiplying the last entry of the preceding row by \( \omega \). In particular, \( \{1\} \)-circulant matrices are circulant matrices, while \( \{-1\} \)-circulant matrices are skew-circulant matrices. Huckle [111] and R. Chan and Jin [35] have used skew-circulant matrices as preconditioners for Toeplitz matrices and proved that under the Wiener class assumptions, the spectra of these preconditioned matrices are clustered around 1. Performances of general \( \{ \omega \} \)-circulant matrices as preconditioners for Toeplitz matrices are discussed in R. Chan and Ng [46] and Huckle [113].

2.5. General remarks on circulant preconditioners. Throughout the preceding sections, we discussed many different kinds of circulant preconditioners. We note from Theorem 2.7 that most of them can be derived from the convolution approach. Moreover, the theorem changes the problem of finding a preconditioner to a problem in approximation theory. In particular, using results in approximation theory, the theorem can give us a guideline as to which preconditioner is better for a given generating function. For example, for 2\( \pi \)-periodic continuous functions that are not in the Wiener class, T. Chan’s preconditioner is better than the Strang preconditioner. Also, if we use a positive kernel to construct the preconditioner, then the preconditioner retains positive definiteness of the given Toeplitz matrix.

We emphasize that the assumptions on the generating functions \( f \) in the theorems in the preceding sections are to simplify the arguments. The main thing required in the proof is not an explicit form of \( f \) but a bound on the rate of decay of the diagonals \( \{ a_j \}_{j=0}^{\infty} \); see the definition of Wiener class functions and also the statement of Theorem 2.4. Thus, for the circulant preconditioning methods to work, there is no need to know the exact form of \( f \), but just an estimate of the decay rate of \( a_j \). Along this line, we remark further that Zygmund [200, p. 183] has shown that if the diagonals \( a_j \) are convex, i.e., the second-order differences \( a_{j+1} - 2a_j + a_{j-1} \geq 0 \) for all \( j \), then \( f \) is nonnegative. Moreover, if one of the second-order differences is positive, then \( f \) is positive.

3. Noncirculant preconditioners for Toeplitz systems.

3.1. Dense Toeplitz preconditioners. As alternatives to circulant preconditioners, Toeplitz matrices have also been proposed and analyzed as preconditioners for Toeplitz systems. We recall that Toeplitz matrix-vector products can be computed in \( O(n \log n) \) operations by using FFTs. Thus, Toeplitz matrices themselves will be good candidates for preconditioners.

In [46], R. Chan and Ng used the Toeplitz matrix \( \tilde{A} \) generated by \( 1/f \) to approximate the inverse of the Toeplitz matrix \( A \) generated by \( f \); i.e., the preconditioned matrix is \( \tilde{A}A \). We remark that the inverse of a Toeplitz matrix is non-Toeplitz in general, but is closely related to Toeplitz matrices; see Friedlander et al. [81]. It has been proven in [46] that the spectrum of the preconditioned matrix \( \tilde{A}A \) is clustered around 1. However, in general it may be difficult to compute the Fourier coefficients of \( 1/f \) explicitly, and hence \( \tilde{A} \) cannot be formed efficiently. R. Chan and Ng [46] thus have derived families of Toeplitz preconditioners \( P^{(s)} \) by using different kernel functions mentioned in §2.2 and different levels of approximation for the Fourier coefficients of \( 1/f \). For the first level of approximation, \( s = 1 \), \( P^{(1)} \) is the circulant preconditioner mentioned in §2.2, depending on which kernel function is used. For integers \( s > 1 \), the preconditioner \( P^{(s)} \) thus constructed can be written as a sum of \( \{ \omega \} \)-circulant matrices. More precisely,

\[
P^{(s)} = \sum_{t=0}^{s-1} W^{(t)},
\]

where \( W^{(t)} \) are \( e^{-2\pi i t/s} \)-circulant matrices.
Using this fact, we showed that given any Toeplitz matrix $A$ and integer $s > 1$, we can decompose $A$ as

$$A = \sum_{t=0}^{s-1} U^{(t)},$$

where $U^{(t)}$ are $\omega$-circulant matrices. For $s = 2$, the formula just states that any Toeplitz matrix can be written as the sum of a circulant matrix and a skew-circulant matrix, a fact first discovered by Pustyl'nikov [161]. Using the decomposition, we showed further that if all $U^{(t)}$ are invertible, then

$$P^{(s)} = \sum_{t=0}^{s-1} (U^{(t)})^{-1}.$$

We recall that in the additive Schwarz method for elliptic problems, a matrix $A$ is first decomposed into the sum of individual projection matrices,

$$A = \sum_{t=0}^{s-1} A^{(t)},$$

and then the generalized inverses of these matrices are added back together to form a preconditioner $P$ for the original matrix $A$; i.e.,

$$P = \sum_{t=0}^{s-1} A^{(t)+}$$

(see Dryja and Widlund [72]). In this respect, our construction of the Toeplitz preconditioner $P^{(s)}$ is very similar to the approach used in the additive Schwarz-type preconditioners.

As for the convergence rate, R. Chan and Ng [46] proved that these preconditioned matrices $P^{(s)}A$ have clustered spectra around 1, and numerical results show that under the Wiener class assumption on the Toeplitz matrices this method converges faster than those preconditioned by circulant preconditioners.

Recall that in each iteration of the preconditioned conjugate gradient method, we require the matrix-vector product $P^{(s)}y$. Since $P^{(s)}$ is a Toeplitz matrix, the product can be computed by using FFTs of size $2n$ in $O(2n \log(2n))$ operations. Since circulant matrix-vector multiplication can be done by using FFTs of size $n$, the cost per iteration of this method is roughly $4/3$ times that required by circulant preconditioned systems.

Another way of looking at these Toeplitz preconditioners is by embedding. For a given $n$-by-$n$ Toeplitz matrix $A$, we embed $A$ into a $sn$-by-$sn$ circulant matrix $C_{sn}$ with its first column given by

$$[C_{sn}]_{k,0} = \begin{cases} a_k, & 0 < k < n, \\ 0, & n \leq k \leq sn - n, \\ a_{k-sn}, & sn - n < k < sn. \end{cases}$$

The Toeplitz preconditioner $P^{(s)}$ is just equal to the leading $n$-by-$n$ principal submatrix of $C_{sn}^{-1}$.

Recently, Hanke and Nagy [98] independently considered this approach of constructing Toeplitz preconditioners for band-Toeplitz matrices with bandwidth $2\beta + 1$. The given band-Toeplitz matrix is first embedded into an $(n + \beta)$-by-$(n + \beta)$ circulant matrix $C_{n+\beta}$ by dragging
down the diagonals to fill the northeast and southwest corners of the expanded matrix. Then the inverse of $C_{n+\beta}$ is partitioned as

$$C^{-1}_{n+\beta} = \begin{bmatrix} B^{(1)} & B^{(2)} \\ B^{(3)} & B^{(4)} \end{bmatrix},$$

where $B^{(1)}$ is a square matrix of size $n$. We note that since $C_{n+\beta}^{-1}$ is circulant, $B^{(1)}$ is a Toeplitz matrix. Hanke and Nagy's preconditioned system is defined to be $B^{(1)}A$.

In another development, Linzer [137] proposed using the Schur complement

$$B = B^{(1)} - B^{(2)}(B^{(4)})^{-1}B^{(3)}$$

rather than $B^{(1)}$ in the preconditioning of band-Toeplitz matrices. We note that $B$ is not a Toeplitz matrix in general. However, since a circulant matrix can be diagonalized by the discrete Fourier transform matrix, it follows that $By$ can be computed by using FFTs of size $(n + \beta)$ and by solving the $\beta$-by-$\beta$ Toeplitz system $B^{(4)}z = v$. Linzer [137] showed that $B$ performs better than $B^{(1)}$ theoretically and numerically. We remark that Schur's complement formula was used by Jain [117], Morf [143], and Bitmead and Anderson [16] to develop direct methods for solving Toeplitz systems in signal processing.

3.2. Optimal transform based preconditioners. From (1.4), we see that circulant matrices are precisely those matrices that can be diagonalized by the discrete Fourier transform, a transform that has a fast algorithm for its computations. However, there are other transforms (for instance, the Hartley transform and the sine and cosine transforms) with fast algorithms; see [188]. It is therefore natural to consider using these fast transforms to construct new classes of preconditioners for solving Toeplitz systems.

Recall that the optimal circulant preconditioners $c(A)$, which are the minimizers of $\|C - A\|_F$ over all circulant matrices $C$, are good preconditioners for Toeplitz systems; see §2.1.2. One may therefore consider preconditioners that are minimizers of $\|Q - A\|_F$ over a set of matrices $Q$ that can be diagonalized by a fast discrete transform matrix $\Psi$. According to the terminology used in T. Chan [57] for optimal circulant preconditioners (where the preconditioner is based on the FFT), we call these minimizers the optimal transform based preconditioners.

Since the Frobenius norm is a unitary-invariant norm, the minimum of $\|Q - A\|_F$ over all $Q$ of the form $Q = \Psi \Delta \Psi^*$, $\Delta$ a diagonal matrix, is attained at $\Psi \Delta \Psi^*$. Here $\Delta$ is a diagonal matrix with diagonal entries

$$\Delta_{j,j} = |\Psi^* A \Psi|_{j,j}, \quad j = 1, \ldots, n. \tag{3.1}$$

For Toeplitz matrices, (3.1) can be computed directly for optimal transform based preconditioners. However, computing $\Delta$ using (3.1) is costly even when the matrix-vector product $\Psi y$ can be done efficiently. We emphasize that to construct T. Chan's minimizers economically, we exploit the fact that the class of circulant matrices has a very nice basis, namely, the shift operator (i.e., the circulant matrix with $[0, \ldots, 0, 1]$ as its first row) and its powers. Given another transform matrix $\Psi$, in order to construct its minimizer efficiently, we need to find matrices having special algebraic structures to characterize all matrices that can be diagonalized by $\Psi$. This is the crucial step in finding a fast algorithm for obtaining the minimizers.

In the following, optimal sine transform based and optimal Hartley transform based preconditioners for symmetric Toeplitz matrices $A$ are considered. The constructions of such preconditioners require $O(n)$ operations for Toeplitz matrices, the same count as that for the
optimal circulant preconditioner \( c(A) \). Similar to T. Chan’s circulant preconditioner, these optimal transform based preconditioners are also defined for arbitrary matrices. In general, the construction of such optimal approximations for any given \( n \)-by-\( n \) matrix is of \( O(n^2) \) operations; see R. Chan, Ng, and Wong [48].

### 3.2.1. Optimal sine transform based preconditioner

The \((j, k)\) entry of the \( n \)-by-\( n \) discrete sine transform matrix \( \Psi^{(s)} \) is given by

\[
\sqrt{\frac{2}{n+1}} \sin \left( \frac{\pi j k}{n+1} \right), \quad 1 \leq j, k \leq n.
\]

For any \( n \)-vector \( v \), the matrix-vector product \( \Psi^{(s)} v \) can be done in \( O(n \log n) \) operations by the fast sine transforms; see, for instance, Yip and Rao [196]. We define the optimal sine transform based preconditioner \( s(A) \) to be the minimizer of \( \| Q - A \|_F \) over the set of matrices \( Q \) that can be diagonalized by \( \Psi^{(s)} \).

Let \( S \) be the vector space over \( \mathbb{R} \) containing all \( n \)-by-\( n \) matrices that can be diagonalized by the discrete sine transform matrix \( \Psi^{(s)} \), i.e.,

\[
S = \{ \Psi^{(s)} \Lambda \Psi^{(s)} \mid \Lambda \text{ is a diagonal matrix} \}.
\]

Boman and Koltracht [19], Bini and Di Benedetto [15], and Huckle [114] independently proved that a matrix belongs to \( S \) if and only if the matrix can be expressed as a special sum of a Toeplitz matrix and a Hankel matrix. The idea of their proof is to exhibit a basis for \( S \) with each element in the basis being a sparse matrix. The following theorem gives the basis that Boman and Koltracht derived.

**Theorem 3.1** (see Boman and Koltracht (1995) [19]). Let \( Z_i, i = 1, \ldots, n \), be \( n \)-by-\( n \) matrices with the \((h, k)\) entry given by

\[
Z_i(h, k) = \begin{cases} 
1 & \text{if } |h - k| = i - 1, \\
-1 & \text{if } h + k = i - 2, \\
-1 & \text{if } h + k = 2n - i + 3, \\
0 & \text{otherwise.}
\end{cases}
\]

Then \( \{ Z_i \}_{i=1}^n \) is a basis for \( S \).

Thus, the vector space \( S \) can be identified as follows.

**Theorem 3.2** (see Boman and Koltracht (1995) [19], Bini and Di Benedetto (1990) [15], and Huckle (1994) [114]). Any matrix \( Q \) in \( S \) can be written as \( Q = X - Y \), where \( X \) is a symmetric Toeplitz matrix with first column \( x = [x_1, x_2, \ldots, x_n]^T \) and \( Y \) is a Hankel matrix with first column \( [0, 0, x_n, \ldots, x_3]^T \) and last column \( [x_3, \ldots, x_n, 0, 0]^T \).

Similar to the case of circulant matrices (cf. (1.5)), there is a relationship between the first column of matrices \( Q \) in \( S \) and their eigenvalues. For any \( Q \) in \( S \), its eigenvalues can be obtained by multiplying the first column of \( Q \) by \( \Psi^{(s)} \), which can be done in \( O(n \log n) \) operations by fast sine transforms. Hence, any matrix in \( S \) is determined by its first column. The following theorem gives the explicit formula for the entries of the first column of \( s(A) \).
THEOREM 3.3 (see R. Chan, Ng, and Wong [48]). Let $A$ be an $n$-by-$n$ symmetric Toeplitz matrix. Let $[s_1, s_2, \ldots, s_n]^T$ be the first column of $s(A)$. Then

$$s_k = \begin{cases} 
\frac{a_0 - \left(\frac{n - 2}{n + 1}\right) a_2}{n + 1}, & k = 1, \\
\frac{a_1 - \left(\frac{n - 3}{n + 1}\right) a_3}{n + 1}, & k = 2, \\
\left(\frac{n - k + 3}{n + 1}\right) a_{k-1} - \left(\frac{n - k - 1}{n + 1}\right) a_{k+1}, & k = 3, 4, \ldots, n - 2, \\
\left(\frac{4}{n + 1}\right) a_{n-2}, & k = n - 1, \\
\left(\frac{3}{n + 1}\right) a_{n-1}, & k = n.
\end{cases}$$

We see that the cost of constructing the minimizer $s(A)$ is $O(n)$ operations.

3.2.2. Optimal Hartley transform based preconditioner. The $(j, k)$ entry of the $n$-by-$n$ discrete Hartley transform matrix $\Psi^{(h)}$ is given by

$$\frac{1}{\sqrt{n}} \cos \left(\frac{2\pi j k}{n}\right) + \frac{1}{\sqrt{n}} \sin \left(\frac{2\pi j k}{n}\right), \quad 0 \leq j, k \leq n - 1.$$ 

The discrete Hartley transform of any $n$-vector can be computed in $O(n \log n)$ real operations. Bini and Favati [14] first characterized the class of matrices that can be diagonalized by the discrete Hartley transform, and showed that the class is the set of matrices that can be expressed as a special sum of a circulant matrix and a Hankel matrix.

THEOREM 3.4 (see Bini and Favati (1993) [14]). Any matrix $Q$ that can be diagonalized by $\Psi^{(h)}$ can be expressed as

$$(3.3) \quad Q = W + XY,$$

where $W$ is a circulant matrix, $Y$ is a skew-circulant matrix with the first entry of its first column being zero, and

$$X = \begin{bmatrix} 
1 & 0 & \cdots & 0 & 0 \\
0 & \ddots & \ddots & \ddots & 1 \\
& \ddots & \ddots & \ddots & \ddots \\
0 & \ddots & \ddots & \ddots & \ddots \\
0 & 1 & 0 & \cdots & 0 
\end{bmatrix}.$$ 

Using Theorem 3.4, Bini and Favati [14] then determined the minimizer $h(A)$ of $\|Q - A\|_F$ over all matrices that can be diagonalized by $\Psi^{(h)}$. As before, the cost of constructing $h(A)$ is $O(n)$ operations and its eigenvalues can be computed in $O(n \log n)$ operations by using fast Hartley transforms.

THEOREM 3.5 (see Bini and Favati (1993) [14]). Let $A$ be a symmetric Toeplitz matrix. Let $w$ and $y$ be the first columns of $W$ and $Y$ of $h(A)$ as defined in (3.3). Then we have

$$w_k = \frac{(n - k)a_k + k a_{n-k}}{n}, \quad k = 0, 1, \ldots, n - 1.$$
and
\[
y_k = \begin{cases} 
0, & k = 0, \\
\frac{a_k - a_{n-k}}{n}, & k = 1, 2, \ldots, n - 1.
\end{cases}
\]

We note that the optimal circulant preconditioner \( c(A) \) is just equal to the circulant part of \( h(A) \). As for how good the optimal sine transform based and Hartley transform based preconditioners are as preconditioners for Toeplitz systems, R. Chan, Ng, and Wong [48], Bini and Favati [14], Bini and Di Benedetto [15], and Jin [119] proved that they both have the same convergence properties as the optimal circulant preconditioners. More precisely, if a given Toeplitz matrix \( A \) is generated by a \( 2\pi \)-periodic positive continuous function, then both the spectra of \( s(A)^{-1}A \) and \( h(A)^{-1}A \) are clustered around 1. However, numerical results given in R. Chan, Ng, and Wong [48], Bini and Favati [14], and Jin [119] show that the convergence performance of these new transform based preconditioners is better in terms of the number of iterations than that of the optimal circulant preconditioner.

In addition, Boman and Koltracht [19] have constructed bases for the vector spaces of matrices that can be diagonalized by some common fast transform matrices in engineering. These transform matrices include different cosine transform matrices and another sine transform matrix defined in [188]. Thus, one can construct optimal transform based preconditioners corresponding to these transform matrices. Boman and Koltracht [19] and R. Chan, Ching, and Wong [33] have studied optimal cosine transform based preconditioners. The cost of construction is the same as those of \( s(A) \) and \( h(A) \).

### 3.3. Band-Toeplitz preconditioners.

In this subsection, we consider Toeplitz matrices \( A \) generated by nonnegative \( 2\pi \)-periodic real-valued functions. We first recall that a function \( f \) is said to have a \( v \)th order zero at \( \theta_0 \) if \( f(\theta_0) = 0 \) and \( v \) is the smallest positive integer such that \( f^{(v)}(\theta_0) \neq 0 \) and \( f^{(v+1)}(\theta) \) is continuous in a neighborhood of \( \theta_0 \). With the knowledge of the order of \( f \) at its minimum, we can give a better estimate of the spectrum of \( A \) than that in (1.2).

**Theorem 3.6** (see R. Chan (1991) [29]). Suppose that \( f(\theta) - f_{\min} \) has a unique zero of order \( 2v \) at \( \theta = \theta_0 \). Then for all \( n > 0 \), we have
\[
\lambda_{\text{min}}(A) \leq d_1 f_{\min} + d_2 n^{-2v},
\]
and
\[
\kappa(A) \geq \frac{d_3 n^{2v}}{d_4 + f_{\min} n^{2v}},
\]
where \( \{d_i\}_{i=1}^4 \) are some constants independent of \( n \).

Thus, when \( f_{\min} = 0 \), the condition number of \( A \) is not uniformly bounded and the Toeplitz matrix \( A \) is ill conditioned. Tyrtyshnikov has proven theoretically [185] that Strang's and T. Chan's preconditioners will fail in this case. In fact, he showed that the numbers of outlying eigenvalues of \( S^{-1}A \) and \( c(A)^{-1}A \) are of \( O(n^{v/(v+\mu)}) \) and \( O(n^{v/(v+1)}) \), respectively. Here, \( \mu \) is the degree of smoothness of the function \( f \) and \( v \) is the order of \( f \) at the zeros. These results were numerically verified in Tyrtyshnikov and Strela [187].

Instead of finding other possible circulant preconditioners, R. Chan [29] resorted to using band-Toeplitz matrices as preconditioners. The motivation behind using band-Toeplitz matrices is to approximate the generating function \( f \) by trigonometric polynomials of fixed degree rather than by convolution products of \( f \) with some kernels. The advantage here is that trigonometric polynomials can be chosen to match the zeros of \( f \), so that the preconditioned method still works when \( f \) has zeros.
THEOREM 3.7 (see R. Chan (1991) [29] and R. Chan and Ng (1993) [45]). Let $f$ be a nonnegative piecewise continuous real-valued function defined on $[-\pi, \pi]$. Suppose that $f(\theta) - f_{\text{min}}$ has a unique zero of order $2v$ at $\theta = \theta_0$. Let $B$ be the Toeplitz matrix generated by the function

$$b_v(\theta) = [2 - 2\cos(\theta - \theta_0)]^v + f_{\text{min}}.$$  

Then $\kappa(B^{-1}A)$ is uniformly bounded for all $n > 0$.

We note that $B$ is a band matrix with bandwidth $2v + 1$ and its diagonals can be obtained by using Pascal's triangle. The band system $By = z$ can be solved by using any band matrix solver; see Golub and Van Loan [88] or Wright [193] for a parallel one. The cost of factorizing $B$ is about $\frac{1}{2}v^2n$ operations, and then each subsequent solve requires an extra $(2v + 1)n$ operations. Hence, the total number of operations per iteration is of $O(n \log n)$, as $v$ is independent of $n$.

When $f_{\text{min}} = 0$, the band preconditioner has improved the condition number from $\kappa(A) = O(n^{2v})$ to $\kappa(B^{-1}A) = O(1)$. Since the number of iterations required to attain a given tolerance $\epsilon$ is bounded by

$$\frac{1}{2} \sqrt{\kappa(B^{-1}A)} \log \left( \frac{2}{\epsilon} \right) + 1$$

(see, for instance, [7, p. 26]), the overall work required to attain the given tolerance is reduced from $O(n^{v+1} \log n)$ to $O(n \log n)$ operations. As for the storage, we just need an $n$-by-$n$ $(2v + 1)$ matrix to hold the factors of the preconditioner $B$. Thus, the overall storage requirement in the conjugate gradient method is about $(8 + v)n$. Finally, we remark that similar results hold when there are multiple points on $[-\pi, \pi]$, where $f$ takes on its minimum value; see R. Chan [29].

The main drawback of using these band-Toeplitz matrices as preconditioners is that when $f$ is positive, these preconditioned systems converge much slower than those preconditioned by circulant preconditioners. Therefore, R. Chan and Tang [50] designed other kinds of band-Toeplitz preconditioners such that their preconditioned systems converge at the same rate as the circulant preconditioned systems even when $f$ is positive. Their idea is to increase the bandwidth of the band-Toeplitz preconditioner to get extra degrees of freedom, which enable them not only to match the zeros in $f$, but also to minimize the relative error $\| (f - g)/f \|_\infty$ in approximating $f$ by trigonometric polynomials $g$. The minimizer, which is a trigonometric polynomial, is found by a version of the Remez algorithm proposed by Tang [177].

THEOREM 3.8 (see R. Chan and Tang (1994) [50]). Let $f$ be the generating function of $A$ and $g_\ell$ be the minimizer of $\| (f - g)/f \|_\infty$ over all trigonometric polynomials of degree $\ell$. If

$$\left\| \frac{f - g_\ell}{f} \right\|_\infty = \alpha < 1,$$

then the Toeplitz matrix $B_\ell$ generated by $g_\ell$ is positive definite and

$$\kappa(B_\ell^{-1}A) \leq \frac{1 + \alpha}{1 - \alpha}, \quad n = 1, 2, 3, \ldots.$$ 

The parameter $\alpha$ is given explicitly in the Remez algorithm. It gives an a priori bound on the number of iterations required for convergence.

The main idea behind Theorems 3.7 and 3.8 is to approximate the given nonnegative generating function $f$ by trigonometric polynomials that match the zeros of $f$. Clearly, any function $g$ that matches the zeros of $f$ and gives rise to Toeplitz matrices that are easily
invertible can also be considered. This idea is exploited in Di Benedetto [9], Di Benedetto, Fiorentino, and Serra [11], and Serra [168]. In [11], \( f \) is first approximated by \( b_\nu \) as in (3.4), and then the quotient \( f/b_\nu \) is further approximated by a trigonometric polynomial or rational function to enhance the convergence rate. In [168], \(|f|\) is used to generate a Toeplitz preconditioner for the Toeplitz matrices generated by nondefinite functions \( f \).

From the above discussion, we see that if \( f_{\text{min}} = 0 \), then the circulant preconditioners fail because they cannot match the zeros of the given \( f \), while the band-Toeplitz preconditioners give only linear convergence because the preconditioned matrices do not have clustered spectra. R. Chan and Ching [32] considered using products of circulant matrices and band-Toeplitz matrices as preconditioners for Toeplitz systems generated by nonnegative functions. The band-Toeplitz part of these \textit{circulant-Toeplitz preconditioners} is to match the zeros of the given function, and the circulant part is to speed up the convergence rate of the algorithm. Instead of using powers of \( 2 - 2 \cos \theta \) as in (3.4) to generate the band-Toeplitz part of the preconditioner, they considered using powers of \( 1 - e^{i\theta} \) instead. This results in preconditioners that can handle complex-valued generating functions with zeros of arbitrary orders. We remark that Freund and Huckle [79] also considered using the band-Toeplitz and circulant preconditioning techniques via a displacement-based formula to solve these Toeplitz systems without knowledge of the underlying generating function.

Another approach for handling ill-conditioned Toeplitz matrices \( A \) is developed in Concus and Saylor [67]. The Trench algorithm [182] is performed on \( A \). When breakdowns or near-breakdowns occur, the matrix is perturbed so that the algorithm can proceed to produce an approximate inverse \( B \). The resulting preconditioned system is \( BA \), and they showed that it has a clustered spectrum.

One application of band-Toeplitz preconditioners is in the solving of Toeplitz-plus-band systems \( (A + D)x = b \). Here \( A \) is an \( n \times n \) Hermitian Toeplitz matrix and \( D \) is an \( n \times n \) Hermitian band matrix with bandwidth independent of \( n \). These systems appear in solving Fredholm integrodifferential equations of the form

\[
D[x(\theta)] + \int_{\alpha}^{\beta} a(\theta - \phi)x(\phi)d\phi = b(\theta),
\]

where \( x(\theta) \) is the unknown function to be found, \( a(\theta) \) is a convolution kernel, and \( D \) is a differential operator. After discretization, the integral will lead to a Toeplitz matrix and \( D \) to a band matrix; see Delves and Mohamed [71, p. 343]. Toeplitz-plus-band matrices also appear in signal processing literature and have been referred to as peripheral innovation matrices; see Carayannis, Kalouptsidis, and Manolakis [24].

Unlike Toeplitz systems, there exist no fast direct solvers for solving Toeplitz-plus-band systems. It is mainly because the displacement rank of the matrix \( A + D \) can take any value between 0 and \( n \). Hence, fast Toeplitz solvers that are based on a small displacement rank of the matrices cannot be applied. Conjugate gradient methods with circulant preconditioners do not work for Toeplitz-plus-band systems either. In fact, Strang's circulant preconditioner is not even defined for non-Toeplitz matrices. T. Chan's circulant preconditioner, while defined for \( A + D \), does not work well when the eigenvalues of \( D \) are not clustered; see [45]. Also, the matrix \( c(A) + D \) cannot be used as a preconditioner, for it cannot be inverted easily. In [45], R. Chan and Ng proposed using the matrix \( B + D \) to precondition \( A + D \), where \( B \) is the band-Toeplitz preconditioner given in Theorem 3.7. Both theoretical and numerical results show that the convergence performance of the preconditioner is better in the number of iterations than those of nonpreconditioned systems and circulant preconditioned systems. Clearly, instead of \( B \), the band-Toeplitz preconditioners in Theorem 3.8 could be used.
4. Applications to Toeplitz-related systems. In the following sections, we discuss applications of optimal transform based preconditioners to Toeplitz-related systems arising from partial differential equations, queueing problems, signal and image processing, integral equations, and time series analysis. Part of the motivation for using optimal transform based preconditioners is to exploit their fast inversion via their transform matrices. In most of the applications, we will simply use the optimal circulant-type preconditioners. We therefore start by extending the results for point optimal circulant preconditioners in §2.1.2 to block-circulant preconditioners.

4.1. Block-circulant preconditioners. Let us consider a general system $Ax = b$, where $A$ is an $mn$-by-$mn$ matrix partitioned as

\[
A = \begin{bmatrix}
A_{1,1} & A_{1,2} & \cdots & A_{1,m} \\
A_{2,1} & A_{2,2} & \cdots & A_{2,m} \\
\vdots & \vdots & \ddots & \vdots \\
A_{m,1} & A_{m,2} & \cdots & A_{m,m} 
\end{bmatrix}.
\]

Here the blocks $A_{i,j}$ are square matrices of order $n$. Given such a matrix $A$, one obvious choice is to use the $mn$-by-$mn$ point-circulant matrix $c(A)$, defined by (2.6), as a circulant approximation to $A$. However, this in general will spoil the block structure of $A$.

T. Chan and Olkin [59] and Holmgren and Otto [105], in solving noise reduction problems and hyperbolic differential equations, independently proposed using circulant-block (CB) matrices (cf. Davis [69, p. 181]) to approximate $A$. Since $c(\cdot)$ is well defined for any square matrix, it is natural to define the CB approximation to $A$ as

\[
c_1(A) = \begin{bmatrix}
c(A_{1,1}) & c(A_{1,2}) & \cdots & c(A_{1,m}) \\
c(A_{2,1}) & c(A_{2,2}) & \cdots & c(A_{2,m}) \\
\vdots & \vdots & \ddots & \vdots \\
c(A_{m,1}) & c(A_{m,2}) & \cdots & c(A_{m,m}) 
\end{bmatrix}.
\]

Some of the spectral properties of $c(A)$ can be extended to $c_1(A)$ (cf. (2.7)).

**Theorem 4.1** (see R. Chan and Jin [1992] [36]). *Given any mn-by-mn Hermitian matrix $A$ partitioned as in (4.1), we have

\[
\lambda_{\min}(A) \leq \lambda_{\min}(c_1(A)) \leq \lambda_{\max}(c_1(A)) \leq \lambda_{\max}(A).
\]

In particular, if $A$ is positive definite, then $c_1(A)$ is also positive definite. Moreover, the operator $c_1(\cdot)$ is a linear projection operator with operator norm $\|c_1\|_2 = \|c_1\|_F = 1$.*

It is interesting to note that the matrix $c_1(A)$ is just the minimizer of $\|A - C\|_F$ over all matrices $C$ that are $m$-by-$m$ block matrices with $n$-by-$n$ circulant blocks. It can also be viewed as the approximation of $A$ along one specific direction. It is natural to consider the preconditioner that results from approximation along the other direction. In this case, the preconditioners are block-circulant (BC) matrices (cf. Davis [69, pp. 176–177]). Both CB and BC preconditioners are called *level-1* preconditioners in T. Chan and Olkin [59].

Clearly one can do approximations in both directions and get a preconditioner $c_2(A)$, which is based on circulant approximations within each block and also on each block level. The resulting preconditioners are block-circulant-circulant-block (BCCB) matrices (cf. Davis [69, p. 184]) and are called *level-2* preconditioners in [59]. BCCB preconditioners for block-Toeplitz-Toeplitz-block (BTTB) matrices (cf. T. Chan and Olkin [59]) and low-rank perturbations thereof have been investigated by Holmgren and Otto [104], Ku and Kuo [129], Tytryshnikov [184], and R. Chan and Jin [36].
BTB matrices occur in many applications. Let us consider the cost of constructing CB and BCCB preconditioners for them. Let

\[
A = \begin{bmatrix}
A_0 & A_1 & \cdots & A_{m-1} \\
A_1 & A_0 & \cdots & A_{m-2} \\
\vdots & \vdots & \ddots & \vdots \\
A_{m-1} & A_{m-2} & \cdots & A_0
\end{bmatrix},
\]

where the blocks \(A_k\) are themselves symmetric Toeplitz matrices of order \(n\). By (4.2), the blocks of \(c_1(A)\) are just \(c(A_k)\). From (1.4) and (3.1), we see that \(c(A_k) = F^*\delta(FA_k F^*)F\), where \(\delta(FA_k F^*)\) is the diagonal matrix whose diagonal equals that of \(FA_k F^*\). Since \(A_k\) are Toeplitz matrices, each \(\delta(FA_k F^*)\) can be computed in \(O(n \log n)\) operations. Therefore, we need \(O(mn \log n)\) operations to form

\[
(4.4) \quad \Delta = (I \otimes F)c_1(A)(I \otimes F^*) = \begin{bmatrix}
\delta(FA_0 F^*) & \delta(FA_1 F^*) & \cdots & \delta(FA_{m-1} F^*) \\
\delta(FA_1 F^*) & \delta(FA_0 F^*) & \cdots & \delta(FA_{m-2} F^*) \\
\vdots & \vdots & \ddots & \vdots \\
\delta(FA_{m-1} F^*) & \delta(FA_{m-2} F^*) & \cdots & \delta(FA_0 F^*)
\end{bmatrix}
\]

Here \(\otimes\) is the Kronecker tensor product and \(I\) is the \(m\)-by-\(m\) identity matrix.

To solve \(c_1(A)y = d\), we permute the diagonals in (4.4) to form a block diagonal matrix

\[
(4.5) \quad \tilde{A} = P^*\Delta P = \begin{bmatrix}
\tilde{A}_{1,1} & 0 & \cdots & 0 \\
0 & \tilde{A}_{2,2} & \cdots & 0 \\
\vdots & \vdots & \ddots & \vdots \\
0 & 0 & \cdots & \tilde{A}_{n,n}
\end{bmatrix},
\]

where

\[
[\tilde{A}_{k,k}]_{ij} = [\delta(FA_{i,j} F^*)]_{kk} = [\delta(FA_{i-j} F^*)]_{kk}, \quad 1 \leq i, j \leq m, \quad 1 \leq k \leq n.
\]

Notice that the diagonal blocks \(\tilde{A}_{k,k}\) are still symmetric Toeplitz matrices of order \(m\). We note that the linear equations with \(\tilde{A}_{k,k}\) must be solved in every iteration step. Therefore, it is convenient to use \(O(m \log^2 m)\) operations algorithm (cf. [4]) to compute a Gohberg–Semencul formula [87]. Thus, the matrix-vector product \(\tilde{A}_{k,k}^{-1}v\) for any vector \(v\) can be computed in \(O(m \log m)\) operations. Hence the system \(c_1(A)y = d\) can be solved in \(O(mn(\log m + \log n))\) operations.

For the BCCB preconditioner \(c_2(A)\), we must take the level-1 approximation of \(\tilde{A}\) in (4.5). That will add another \(O(nm \log m)\) operations to the construction cost. However, the cost of solving \(c_2(A)y = d\) is reduced from \(O(nm \log^2 m + mn(\log m + \log n))\) to \(O(mn \log(mn))\) since the preconditioner \(c_2(A)\) can be completely diagonalized by FFTs.

Next we consider the cost of the matrix-vector multiplication \(Av\). We recall that the matrix-vector multiplication \(A_k w\) for any \(n\)-vector \(w\) can be computed by FFTs of size \(2n\) by first embedding \(A_k\) into a \(2n\)-by-\(2n\) circulant matrix; see (1.6). For the matrix-vector product \(Av\), we use the same trick. We first embed \(A\) into a (blockwise) \(2m\)-by-\(2m\) BC matrix, where each block itself is a \(2n\)-by-\(2n\) circulant matrix. Then we extend \(v\) to a \(4mn\)-vector by putting zeros in the appropriate places. We note that \(Av\) can be obtained in \(O(mn \log(mn))\) operations.

Thus, we conclude that for level-1 circulant preconditioners \(c_1(A)\), the initialization cost is \(O(mn \log n)\) and the cost per iteration is \(O(mn \log^2 m + mn \log n)\). For level-2 preconditioners
\( c_2(A) \), both costs are \( O(mn \log(mn)) \). Hence, it is cheaper to use level-2 preconditioners. As for the convergence rate, we have the following theorem for \( c_2(A) \) when \( A \) is generated by a positive Wiener class function. (See Ku and Kuo [129] for matrices generated by rational functions instead.)

**Theorem 4.2** (see Tyryshnikov [184]). Let \( A \) be given by (4.3) with the entries of the block \( A_j \) denoted by \( a_{pq}^{(j)} = a_{p-q}^{(j)} \) for \( 0 \leq p, q < n, 0 \leq j < m \). If \( a_k^{(j)} \) are Fourier coefficients of a positive function in the Wiener class, i.e.,

\[
\sum_{j=0}^{\infty} \sum_{k=0}^{\infty} |a_k^{(j)}| < \infty
\]

and

\[
\sum_{j=0}^{\infty} \sum_{k=0}^{\infty} a_k^{(j)} e^{-ijx} e^{-iky} > 0 \quad \forall x, y \in [-\pi, \pi],
\]

then the spectrum of \( c_2(A)^{-1} A \) is clustered around 1 for large \( n \) and \( m \).

As a consequence, when the conjugate gradient method is applied to solve the system \( c_2(A)^{-1} A x = c_2(A)^{-1} b \), we expect superlinear convergence as in the point case. We recall that the algorithm requires \( O(mn \log(mn)) \) operations both in the initialization step and also in each iteration. Thus, the total complexity of the algorithm is bounded above by \( O(mn \log(mn)) \).

Besides extending results in the point-circulant case to the BC case, one can also extend the results for other types of preconditioners as discussed in §3 to the block case. Works in this direction can be found in Di Benedetto [10] where sine transform-type preconditioners are extended, in Serra [167] where Theorem 3.7 is extended, in Jin [120] where Theorem 3.8 is extended, in Hemmingsson [100] where Toeplitz-block (TB) and BTTB matrices as preconditioners are considered, and finally in Tyryshnikov [183–186] where results on superoptimal circulant preconditioners and product preconditioners are generalized to the block case.

### 4.2. Applications to partial differential equations.

In this subsection, we consider using preconditioned conjugate gradient methods with optimal transform based preconditioners to solve different types of partial differential equations.

#### 4.2.1. Elliptic problems.

Consider the elliptic problem

\[
(-a(x, y)u_x)_x - (b(x, y)u_y)_y = g(x, y)
\]

on the unit square \([0, 1] \times [0, 1]\) with Dirichlet boundary conditions. After discretization with a mesh size of \(1/(n + 1)\), such a problem reduces to the solution of an \(n^2\)-by-\(n^2\) linear system of the form \( A x = b \). With a standard 5-point stencil and lexicographical ordering, the discretization matrix \( A \) is a block tridiagonal matrix, where the diagonal blocks are tridiagonal matrices and the off-diagonal blocks are diagonal matrices. The linear system is often solved by iterative methods such as the preconditioned conjugate gradient method.

R. Chan and T. Chan [31] proposed two choices of circulant preconditioners for these discretization matrices \( A \). The first one is \( C_P = c(A) + \rho n^{-2} I \), where \( c(\cdot) \) is defined by (2.6). The diagonals of this point-circulant preconditioner are thus obtained as the simple averages of the coefficients \( a(x, y) \) and \( b(x, y) \) over the whole grid. The second choice is a BCCB preconditioner which preserves the block structure of \( A \) and is defined as \( C_B = c_2(A) + \rho n^{-2} I \), where \( c_2(A) \) is the level-2 preconditioner of \( A \) as defined in §4.1. Thus, the diagonals are obtained as the simple averages of the coefficients along the lines of the grid. The constant \( \rho n^{-2} \) added to the main diagonal is to minimize the condition number of the resulting preconditioned
systems, an approach similar in idea to that used in modified incomplete LU factorizations; see [73].

We note that the product $C_B^{-1}y$ can be computed by using $2n$ FFTs of size $n$ and the
solution of $n$ intermediate circulant tridiagonal systems, each requiring $O(n)$ complexity,
whereas $C_P^{-1}y$ requires two FFTs of size $n^2$. Similar circulant preconditioners can be defined
for more general elliptic operators with more complicated difference stencils and also in higher
dimensions.

THEOREM 4.3 (see R. Chan and T. Chan (1992) [31]). Assume that in (4.6)

$$0 < c_{\min} \leq a(x, y), b(x, y) \leq c_{\max}$$

for some constants $c_{\min}$ and $c_{\max}$. Then we have

$$\kappa(C_B^{-1}A) \leq O(n) \quad \text{and} \quad \kappa(C_P^{-1}A) \leq O(n \log n).$$

We remark that for second-order elliptic problems, the condition number $\kappa(A)$ of $A$ is
$O(n^2)$. It follows from the above theorem that the condition number of the system is reduced
from $O(n^2)$ to $O(n)$ by circulant preconditioning. Works in this direction can be found in
Huckle [115], where skew-circulant preconditioners are used for these discretization matrices.
Finally, we note that the application of the circulant preconditioners requires $O(n^2 \log n)$
flops per iteration, which is slightly more expensive than the $O(n^2)$ flops for the incomplete
LU-type preconditioners studied in [7, 73, 94]. However, FFTs can be computed in $O(\log n)$
parallel steps with $O(n^2)$ processors, whereas the incomplete LU preconditioners require at
least $O(n)$ steps regardless of how many processors are available. The computation (based
on averaging the coefficients of the elliptic operator) of these circulant preconditioners is also
highly parallelizable across a wide variety of architectures.

One can view the circulant preconditioners as approximations to the original given equation,
but with the given boundary conditions replaced by periodic ones. It is thus natural
to consider using other fast transform based preconditioners to precondition elliptic
problems. The discrete sine transform matrix $\Psi^{(s)}$ with its entries defined by (3.2) diagonalizes
all symmetric tridiagonal Toeplitz matrices, in particular the 1-d discrete Laplacian with Dirichlet
boundary conditions—tridiag$[-1, 2, -1]$. Therefore, one expects that the optimal sine
transform based preconditioners discussed in §3.2.1 will give better approximations to elliptic
problems with Dirichlet boundary conditions.

In [52], R. Chan and Wong proposed using these optimal sine transform based preconditioners
for matrices $A$ that come from the discretization of second-order elliptic operators.
For simplicity, we let $(\Gamma + \Sigma)\Sigma^{-1}(\Gamma^* + \Sigma)$ be the block Cholesky factorization of $A$
with lower block triangular matrix $\Gamma$ and diagonal block matrix $\Sigma$. For such factorizations,
the preconditioner is defined to be the matrix

$$P = (\hat{\Gamma} + \Phi)\Phi^{-1}(\hat{\Gamma}^* + \Phi)$$

with block diagonal matrix $\Phi$ and lower block triangular matrix $\hat{\Gamma}$. Here the diagonal blocks
of $\Phi$ and the subdiagonal blocks of $\hat{\Gamma}$ are, respectively, the optimal sine transform approximations
to the diagonal blocks of $\Sigma$ and the subdiagonal blocks of $\Gamma$. R. Chan and Wong showed that for 2-d domains, the construction cost of $P$ and the cost for each iteration of
the preconditioned conjugate gradient algorithm are of $O(n^2 \log n)$. They also showed that
for rectangular domains, $P$ can be obtained from $A$ by taking the optimal sine transform approximations of each subblock of $A$. Thus, the construction of $P$ is similar to the level-1
circulant preconditioners for rectangular domains, except that the construction process can
now be extended to irregular regions as well.
TABLE 4.1

Number of iterations for the unit square.

<table>
<thead>
<tr>
<th>$1/h$</th>
<th>$I$</th>
<th>$C_B$</th>
<th>$P$</th>
<th>MILU</th>
<th>MINV</th>
<th>$I$</th>
<th>$C_B$</th>
<th>$P$</th>
<th>MILU</th>
<th>MINV</th>
</tr>
</thead>
<tbody>
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<td>10</td>
<td>5</td>
<td>6</td>
<td>3</td>
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<td>10</td>
<td>5</td>
<td>6</td>
<td>3</td>
</tr>
<tr>
<td>8</td>
<td>25</td>
<td>12</td>
<td>3</td>
<td>9</td>
<td>5</td>
<td>29</td>
<td>13</td>
<td>7</td>
<td>9</td>
<td>4</td>
</tr>
<tr>
<td>16</td>
<td>24</td>
<td>15</td>
<td>3</td>
<td>13</td>
<td>7</td>
<td>54</td>
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<td>9</td>
<td>14</td>
<td>6</td>
</tr>
<tr>
<td>32</td>
<td>90</td>
<td>20</td>
<td>3</td>
<td>20</td>
<td>11</td>
<td>107</td>
<td>25</td>
<td>11</td>
<td>20</td>
<td>10</td>
</tr>
<tr>
<td>64</td>
<td>186</td>
<td>25</td>
<td>3</td>
<td>28</td>
<td>16</td>
<td>209</td>
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<td>50</td>
<td>13</td>
<td>41</td>
<td>22</td>
</tr>
</tbody>
</table>

TABLE 4.2

Number of iterations for the L-shaped domain.

<table>
<thead>
<tr>
<th>$1/h$</th>
<th>$I$</th>
<th>$P$</th>
<th>MILU</th>
<th>MINV</th>
<th>$I$</th>
<th>$P$</th>
<th>MILU</th>
<th>MINV</th>
</tr>
</thead>
<tbody>
<tr>
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<td>3</td>
<td>9</td>
<td>4</td>
<td>24</td>
<td>7</td>
<td>9</td>
<td>4</td>
</tr>
<tr>
<td>16</td>
<td>40</td>
<td>3</td>
<td>12</td>
<td>6</td>
<td>45</td>
<td>9</td>
<td>13</td>
<td>6</td>
</tr>
<tr>
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</tr>
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<td>64</td>
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<td>25</td>
<td>14</td>
<td>169</td>
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<td>26</td>
<td>12</td>
</tr>
<tr>
<td>128</td>
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<td>21</td>
<td>338</td>
<td>14</td>
<td>37</td>
<td>19</td>
</tr>
</tbody>
</table>

For rectangular regions, the condition number of the preconditioned system $P^{-1}A$ is proven to be of $O(1)$. In contrast, the system preconditioned by the MILU, MINV, and optimal circulant preconditioners is of $O(n)$. We remark that a similar construction of optimal circulant approximations on L-shaped domains has recently been considered by Lirkov and Margenov [138].

In the following, we compare the performance of preconditioned conjugate gradient methods with the optimal circulant preconditioners, optimal sine transform based preconditioners, and MILU- and MINV-type preconditioners. The equation we used is

$$
(4.7) \quad \frac{\partial}{\partial x} \left[ (1 + \varepsilon e^{x+y}) \frac{\partial u}{\partial x} \right] + \frac{\partial}{\partial y} \left[ \left( 1 + \varepsilon^2 \sin(2\pi(x+y)) \right) \frac{\partial u}{\partial y} \right] = g(x, y).
$$

The parameter $\varepsilon$ controls the variation of the coefficient functions. We remark that for this equation there are alternate competitive methods, such as multigrid methods.

The initial guess and the right-hand side are chosen to be random vectors and are the same for all methods. The iteration stops when the residual vector $r_q$ at the $q$th iteration satisfies $\|r_q\|_2/\|r_0\|_2 < 10^{-6}$. All computations were done with double precision on a VAX 6420.

We remark again that the costs per iteration of sine transform based or circulant preconditioners are $O(\log n)$ times more than those required by MINV- or MILU-type preconditioners. However, most of the computations involving fast transform based preconditioners can be done in parallel. Thus, it is difficult to just compare either the time or flop counts. For simplicity, we here compare only the iteration numbers. Tables 4.1 and 4.2 show the numbers of iterations required for convergence for equation (4.7) on the unit square and the L-shaped domain $[0, 1]^2 \setminus [1/2, 1] \times [0, 1/2]$, respectively; see [31, 52]. In the tables, the notation $I$ means that no preconditioner is used and the parameter $h$ is the mesh size. These results show that the optimal sine transform based preconditioner requires fewer iterations than the optimal circulant, MILU, and MINV preconditioners for small $\varepsilon$, i.e., when the variation of the coefficients is small.

4.2.2. Domain decomposition. Let $\Omega \subset \mathbb{R}^2$ be a bounded domain with a piecewise smooth boundary $\partial \Omega$. Given $f \in L^2(\Omega)$, we are interested in finding the weak solution
$u \in H^1(\Omega)$ of the following Neumann problem:

$$a(u, v) = \int_{\Omega} f v dx \quad \forall v \in H^1(\Omega),$$

where

$$a(u, v) = \int_{\Omega} \left( \sum_{i,j=1}^2 a_{ij}(x) \frac{\partial u}{\partial x_i} \frac{\partial v}{\partial x_j} + c(x) uv \right) dx,$$

and $a_{ij}(x), c(x) \in L^\infty(\Omega)$. We assume that the bilinear form $a(u, v)$ is symmetric and there exist positive constants $\alpha_1, \alpha_2$ such that

$$\alpha_1 \|v\|^2_{H^1(\Omega)} \leq a(u, v) \leq \alpha_2 \|v\|^2_{H^1(\Omega)}.$$

Suppose that $\Omega$ is a disjoint union of two subdomains $\Omega_1$ and $\Omega_2$ and that a fast solver is available on each subdomain. Domain decomposition methods for elliptic problems defined on a union of subdomains have been studied in many papers; see, for instance, [72, 191]. The idea of substructuring is to reduce the problem in $\Omega$ to one on the interface $\partial \Omega_1 \cap \partial \Omega_2$. Let $H^{1/2}(\partial \Omega_k)$ be the Sobolev space of order one-half on the boundary $\partial \Omega_k$, $k = 1, 2$, with weighted norm:

$$\|v\|_{H^{1/2}(\partial \Omega_k)}^2 = \int_{\partial \Omega_k} \int_{\partial \Omega_k} \frac{|v(x(s), y(s)) - v(x(r), y(r))|^2}{\|x(s), y(s)\| - (x(r), y(r))|^2} ds dr + \frac{1}{d_k} \int_{\partial \Omega_k} |v(x(s), y(s))|^2 ds,$$

(4.8)

where $\| \cdot \|$ is the Euclidean distance in $\mathbb{R}^2$. In [126, 127], Kiss and Molnárka proposed using circulant matrices as preconditioners for these elliptic problems. Their idea is to approximate the Euclidean norm $\| \cdot \|$ in (4.8) by

$$|s - r|_k = \min\{|s - r|, \text{length}(\partial \Omega_k) - |s - r|\}.$$

It turns out that the matrix representation $C$ of the resulting approximated bilinear form is circulant, and for the Neumann problem it is a good approximation.

**Theorem 4.4** (see Kiss and Molnárka [1991, 1992] [127, 126]). *For the Neumann problem, the original bilinear form and the approximated one are spectrally equivalent, and hence $\kappa(C^{-1}A) = O(1)$.*

We remark that Dirichlet and mixed boundary problems have also been considered in [126, 127].

**4.2.3. Hyperbolic and parabolic problems.** The idea of circulant preconditioners has also been applied to systems arising from implicit time-marching methods for first-order hyperbolic equations of the form

$$\frac{\partial u}{\partial t} + a(x, y) \frac{\partial u}{\partial x} + b(x, y) \frac{\partial u}{\partial y} = g(x, y);$$

see Holmgren [103], Holmgren and Otto [104–107], and Otto [156]. For such problems, the discretization matrix $A$ is nonsymmetric and often highly nondiagonally dominant, and hence many classical preconditioning techniques are not effective (and sometimes not well defined). For these problems, the circulant preconditioners are often the only ones that work.
Table 4.3
Number of CGS iterations when $\kappa = 100$. Reprinted with permission from S. Holmgren and K. Otto.

<table>
<thead>
<tr>
<th>$1/h$</th>
<th>None</th>
<th>MILU</th>
<th>Block MILU</th>
<th>BCCB</th>
</tr>
</thead>
<tbody>
<tr>
<td>8</td>
<td>48</td>
<td>18</td>
<td>198</td>
<td>17</td>
</tr>
<tr>
<td>16</td>
<td>223</td>
<td>32</td>
<td>&gt;500</td>
<td>20</td>
</tr>
<tr>
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</tr>
<tr>
<td>64</td>
<td>423</td>
<td>69</td>
<td>&gt;500</td>
<td>20</td>
</tr>
<tr>
<td>128</td>
<td>473</td>
<td>89</td>
<td>&gt;500</td>
<td>20</td>
</tr>
<tr>
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<td>135</td>
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</tr>
<tr>
<td>512</td>
<td>&gt;500</td>
<td>&gt;500</td>
<td>&gt;500</td>
<td>19</td>
</tr>
</tbody>
</table>

In [104], comparisons of circulant and BCCB preconditioners with incomplete LU (ILU) and block ILU are given. Some of the circulant-type preconditioners are obtained by changing the boundary conditions from Dirichlet type to periodic ones. It is found that if $\kappa$, the ratio between the time step and the spatial mesh size, is large, then ILU-type preconditioners fail to converge, while the rate of convergence for circulant-type preconditioners remains constant independent of the mesh size. In [105], a framework of CB and BCCB preconditioners is developed, and an analysis of the convergence rate for a BCCB preconditioner is performed for problems that have both periodic and Dirichlet boundary conditions.

An analysis of a CB preconditioner for problems with only Dirichlet boundary conditions is done in [156], which also contains a Fourier analysis of ILU, MILU, and block MILU preconditioned systems. The analysis indicates that the condition number of the MILU preconditioned system remains $O(\kappa)$, whereas the convergence rate of the CB preconditioned systems remains independent of $\kappa$ and mesh size $h$. The convergence analysis was later extended to include a time-independent equation with a weak artificial viscosity; see [106]. Applications of circulant-type preconditioners to the computation of flow in a driven cavity governed by the Navier–Stokes equations can be found in [107]. Implementation of the method on shared memory vector computers and distributed memory computers are studied in [103].

Table 4.3, taken from [104, Table 8], shows the effectiveness of using BCCB matrices as preconditioners for first-order hyperbolic systems when $\kappa$ is large.

In [100], Hemmingsson considered using TB and BTTB preconditioners for the same first-order hyperbolic PDE. The Toeplitz blocks are obtained by minimization as is done in (2.3), except that the minimization is now taken over a restricted set of Toeplitz matrices. The convergence rate of the resulting preconditioned methods is also favorable and is thoroughly analyzed in [101].

Circulant preconditioners for second-order hyperbolic equations have been considered by Jin and R. Chan [121]. In [121], the equation

$$u_{tt} = (a(x, y)u_x)_x + (b(x, y)u_y)_y + g(x, y)$$

is considered. The coefficient matrix, obtained from an implicit time-marching scheme, has a condition number of $O(\kappa^2) + O(h^{-2})$. With circulant-type preconditioners, Jin and R. Chan [121] proved that the condition number is reduced to $O(\kappa) + O(h^{-1})$. The same idea of proof can also be applied to parabolic equations of the form

$$u_t = (a(x, y)u_x)_x + (b(x, y)u_y)_y + g(x, y);$$

see [118]. In [107], Holmgren and Otto investigated circulant-type solvers for mixed hyperbolic–parabolic equations. The second-order terms were considered small or used as an artificial viscosity, such as for the discretized Euler equations in computational fluid dynamics.
4.3. Applications to queueing problems.

4.3.1. Overflow queueing networks. Consider a 2-queue Markovian network with overflow permitted only from queue 1 to queue 2 when queue 1 is full; see [124]. We are interested in finding the steady-state probability distribution vector of the network. Let $\lambda_i$, $\mu_i$, $n_i$, and $s_i$ be, respectively, the input rate, output rate of a single server, buffer size, and number of servers for queue $i$. If the traffic density, defined as $\lambda_i/(s_i\mu_i)$, is close to 1, i.e.,

$$\frac{\lambda_i}{s_i\mu_i} = 1 + O(n_i^{-\alpha})$$

for some $\alpha > 0$, then the queueing problem resembles a second-order elliptic equation on a rectangle with an oblique boundary condition on one side (the side with overflow) and Neumann boundary conditions on the others; see [26].

The SOR method is one of the standard methods for solving this problem; see [124]. However, in [26], the preconditioned conjugate gradient method was also considered, with the preconditioner constructed by changing the oblique boundary condition to the Neumann boundary condition. This preconditioner will be referred to as the “Neumann” preconditioner below. The convergence rate of the preconditioned conjugate gradient method with the Neumann preconditioner is much better than the SOR method, ranging from 10 times faster for small $n_i$ to about 100 times faster for $n_i = 128$. However, inversion of Neumann preconditioners is expensive when the number of servers $s_i > 1$; see [26].

Since the 1-d discrete Laplacian with Neumann boundary conditions can be diagonalized by the cosine transform matrix, one is naturally led to consider optimal cosine transform based preconditioners mentioned in §3.2 for such queueing systems. By using FFTs, each inversion of the preconditioner requires $O(n_1n_2 \log(n_1n_2))$ operations independent of $s_i$. Thus, the cost per iteration is comparable to the $O(n_1n_2)$ operations required by the SOR method.

Numerical results in R. Chan, Ching, and Wong [33] show that the optimal cosine transform based preconditioners performs even better than the Neumann preconditioner. Table 4.4 gives the number of iterations required for convergence when the tolerance is $10^{-6}$. In the test, the parameters $\lambda_i$, $i = 1, 2$, are set to 1 and $\mu_i$ are computed according to (4.9) with $\alpha = 1$. For the SOR method, the optimal relaxation factor, obtained numerically to four significant digits, is used. The results show that the new method is about 10 to 200 times faster than the SOR method.

4.3.2. Queueing networks with batch arrivals. Queueing systems with batch arrivals occur in many applications, such as the telecommunication networks [154] and the loading dock models [166]. Again, we are interested in finding the stationary distribution vector of the network. It will be the normalized null-vector of the generator matrix of the problem.

For this problem, the generator matrix can be written as a sum of a Toeplitz matrix and a rank $s$ matrix, where $s$ is the number of servers. Since the generator matrix is singular, the Toeplitz matrix will have a generating function $f$ that has zeros. R. Chan and Ching [32]
considered using circulant-Toeplitz preconditioners discussed in §3.3 to precondition such systems. The singularity of the generator matrix is canceled by the band-Toeplitz part of the preconditioner, and the circulant part of the preconditioner is used to speed up the convergence of the algorithm. They proved superlinear convergence of the method when \( s \) is independent of \( n \), the size of the queue. Numerical results in [32] verified the fast convergence for small \( s \) and show, moreover, that the convergence is still linear when \( s = n \).

4.4. Applications to signal and image restoration. Image restoration refers to the removal or reduction of degradations (or blur) in an image using a priori knowledge about the degradation phenomena. Applications of image restoration can be found in remote sensing, where details about the photographed terrain must be resolved; in medical imaging, where the diagnosis is based on the clarity of the x-ray radiographs taken; and in space exploration, where images transmitted back to earth by spacecrafts are analyzed. In ground-based imaging, astronomers seek to remove the degradation of astronomical images caused by atmospheric turbulence, which is in part due to the mixing of the warm and cold air layers. When the quality of the images is degraded by blurring and noise, important information remains hidden and cannot be directly interpreted without numerical processing. Our presentation of the image restoration problem here will be brief, and we refer interested readers to Andrews and Hunt [5] and Jain [116] for more detailed discussions.

We begin with a mathematical model of the image restoration problem. The image of an object can be modeled as

\[
 g(\xi, \delta) = \int_{-\infty}^{\infty} \int_{-\infty}^{\infty} a(\xi, \delta; \alpha, \beta)h(\alpha, \beta)d\alpha d\beta + \eta(\xi, \delta),
\]

where \( g(\xi, \delta) \) is the recorded (or degraded) image, \( h(\alpha, \beta) \) is the original image, and the vector \( \eta(\xi, \delta) \) represents an additive noise. The function \( a(\xi, \delta; \alpha, \beta) \) is called the point-spread function (PSF) and represents the degradation of the image. In digital implementation, (4.10) is discretized to obtain the discrete scalar model

\[
 g(i, j) = \sum_{k=1}^{n} \sum_{\ell=1}^{n} a(i, j; k, \ell)h(k, \ell) + \eta(i, j).
\]

In matrix-vector notation, we obtain the linear algebraic form of the image restoration problem,

\[
 \mathbf{g} = A \mathbf{h} + \mathbf{n},
\]

where \( \mathbf{g} \), \( \mathbf{h} \), and \( \mathbf{n} \) are \( n^2 \)-vectors and \( A \) is an \( n^2 \)-by-\( n^2 \) matrix. This is the square image formulation. Often the discretization is chosen so that \( \mathbf{g} \) is a longer vector than \( \mathbf{h} \). In this case, \( A \) is a rectangular \( m^2 \)-by-\( n^2 \) matrix with \( m > n \). The image restoration problem can be stated as follows. Given the observed image \( \mathbf{g} \), the matrix \( A \) which represents the degradation, and possibly the statistics of the noise vector \( \mathbf{n} \), compute an approximation to the original signal \( \mathbf{h} \).

Writing the PSF as \( a(\xi, \delta; \alpha, \beta) \) provides the most general description of the imaging system. This representation allows the PSF to vary with position in both the image and object planes. In this case the PSF is said to be spatially variant, and the matrix \( A \) in (4.11) will have no special structure. Thus, computing a solution to (4.11) can be very expensive. In many practical applications, though, the PSF is spatially invariant; i.e., it acts uniformly across the image and object planes. In particular, it can be written as

\[
 a(\xi, \delta; \alpha, \beta) = a(\xi - \alpha, \delta - \beta),
\]

and the matrix \( A \) it generates is a BTBB matrix.
Because of the ill-conditioning of \( A \), naively solving \( Ah = g \) will lead to extreme instability with respect to perturbations in \( g \); see [5]. Thus, one cannot consider the noise vector \( n \) insignificant in solving the problem. The method of regularization can be used to achieve stability for these problems [1, 17, 162]. In the classical Tikhonov regularization [92], stability is attained by introducing a stabilizing operator \( D \) (called a regularization operator), which restricts the set of admissible solutions. Since this causes the regularized solution to be biased, a scalar \( \mu \), called a regularization parameter, is introduced to control the degree of bias. More specifically, the regularized solution is computed as

\[
\min \left\| \begin{bmatrix} g \\ 0 \end{bmatrix} - \begin{bmatrix} A \\ \mu D \end{bmatrix} h(\mu) \right\|_2.
\]

The term \( D h \) is added in order to smooth the solution \( h \). Choosing \( D \) as a \( k \)th order difference operator matrix forces the solution to have a small \( k \)th order derivative. Notice that if \( D \) is a Toeplitz matrix, then (4.12) reduces to a block-Toeplitz least squares problem.

Thus in general, let us consider the least squares problem

\[
\min_x \| b - Ax \|_2,
\]

where \( A \) is a rectangular Toeplitz or block-Toeplitz matrix. Besides image restoration problems, Toeplitz least squares problems also arise in other important areas, such as the denoising problems; see T. Chan and Olkin [59]. Considerable effort has been devoted to developing fast algorithms for them and most works have been focused on direct methods, such as the fast QR factorization algorithms of Bojanczyk, Brent, and de Hoog [18], Chun and Kailath [63], Cybenko [68], and Sweet [176]. Here we will consider using the preconditioned conjugate gradient method.

Although the classical conjugate gradient algorithm applies only to square Hermitian positive-definite systems, one can still use it to find the solution to (4.13) by applying it to the normal equations in factored form:

\[
A^*(b - Ax) = 0.
\]

The method can be applied without explicitly forming the normal equations matrix \( A^*A \); see Björck [17]. As in the square case, we can precondition the equation to speed up the convergence. Given a nonsingular matrix \( C \), one can use the conjugate gradient method to solve

\[
\min \| b - AC^{-1}y \|_2,
\]

and then set \( x = C^{-1}y \). The cost per iteration of the preconditioned conjugate gradient method is dominated by matrix-vector multiplications with \( A \) and \( A^* \) and by linear system solves with \( C \) as the coefficient matrix. If \( A \) is an \( m \)-by-\( n \) Toeplitz matrix, then matrix-vector multiplications with \( A \) and \( A^* \) can be accomplished in \( O(m \log n) \) operations using FFTs. As discussed in §1.3, the preconditioner matrix \( C \) should be chosen such that the singular values of \( AC^{-1} \) are clustered around 1 and the linear system with coefficient matrix \( C \) can be easily solved. In the following subsections, we survey some of the possible preconditioners.

### 4.4.1. Block-based preconditioners.

We start with the 1-d Toeplitz least squares problems. For the purpose of constructing the preconditioner, we extend the Toeplitz structure of the matrix \( A \) in (4.13) by padding zeros to the bottom left-hand side. In doing so, we may assume without loss of generality that \( m = kn \) for some positive integer \( k \). This padding is only for convenience in constructing the preconditioner and does not alter the original least squares
problem. In the material to follow, we consider the case where $k$ is a constant independent of $n$. More precisely, we consider $kn$-by-$n$ matrices $A$ of the form

$$A = \begin{bmatrix}
A_1 \\
A_2 \\
\vdots \\
A_k
\end{bmatrix},$$

where each square block $A_j$ is a Toeplitz matrix. Notice that if $A$ itself is a rectangular Toeplitz matrix, then each block $A_j$ is necessarily Toeplitz. Also, the matrix in (4.12) is of this form if the regularization operator $D$ is a Toeplitz matrix.

Following [144, 42], for each block $A_j$, the optimal circulant approximation $c(A_j)$ is constructed. Then the block-based preconditioner is defined to be the square circulant matrix $C$ such that

$$C^*C = \sum_{j=1}^{k} c(A_j)^* c(A_j).$$

Notice that each $c(A_j)$ is an $n$-by-$n$ circulant matrix. Hence, they can all be diagonalized by the $n$-by-$n$ discrete Fourier matrix $F$; i.e., $c(A_j) = F^* \Lambda_j F$, where $\Lambda_j$ is diagonal. Thus, the spectrum of $c(A_j)$, $j = 1, \ldots, k$, can be computed in $O(n \log n)$ operations by using FFT and we have

$$C^*C = F^* \sum_{j=1}^{k} (\Lambda_j^* \Lambda_j) F.$$

Clearly $C^*C$ is a circulant matrix and its spectrum can be computed in $O(kn \log n)$ operations. The preconditioner is then given by

$$C = F^* \left( \sum_{j=1}^{k} \Lambda_j^* \Lambda_j \right)^{\frac{1}{2}} F.$$

Recall that the matrix-vector multiplications with $A$ and $A^*$ can be done in $O(m \log n)$ operations by FFT; the cost per iteration in the preconditioned conjugate gradient method is therefore of $O(m \log n)$. The convergence rate of the method depends on the distribution of the singular values of the matrix $AC^{-1}$, which are the same as the square roots of the eigenvalues of the matrix $(C^*C)^{-1}(A^*A)$.

**Theorem 4.5** (see R. Chan, Nagy, and Plemmons (1994) [42]). Suppose the generating functions of the blocks $A_j$ are $2\pi$-periodic continuous functions. If one of these functions has no zeros, then the spectrum of $(C^*C)^{-1}(A^*A)$ is clustered around 1 for sufficiently large $n$.

Thus, if the condition number $\kappa(A)$ of $A$ is of $O(n^\alpha)$, then for sufficiently large $n$ the number of iterations required for convergence is at most $O(\alpha \log n)$ when $\alpha > 0$. Since the number of operations per iteration in the conjugate gradient method is of $O(m \log n)$, the total complexity of the algorithm is therefore of $O(am \log^2 n)$. In the case when $\alpha = 0$, i.e., $A$ is well conditioned, the method converges in $O(1)$ steps. Hence, the complexity is reduced to just $O(m \log n)$ operations.
Next we consider the 2-d Toeplitz least squares problems. Let \( A \) be a \( kmn \)-by-\( mn \) matrix of the form

\[
A = \begin{bmatrix}
A^{(1)} \\
A^{(2)} \\
\vdots \\
A^{(k)}
\end{bmatrix},
\]

(4.14)

where each block \( A^{(i)} \), \( i = 1, \ldots, k \), is a TB matrix. More precisely, \( A^{(i)} \) can be partitioned as

\[
A^{(i)} = \begin{bmatrix}
A_{1,1}^{(i)} & A_{1,2}^{(i)} & \cdots & A_{1,m}^{(i)} \\
A_{2,1}^{(i)} & A_{2,2}^{(i)} & \cdots & A_{2,m}^{(i)} \\
\vdots & \vdots & \ddots & \vdots \\
A_{m,1}^{(i)} & A_{m,2}^{(i)} & \cdots & A_{m,m}^{(i)}
\end{bmatrix}, \quad i = 1, \ldots, k,
\]

where each \( A_{\alpha,\beta}^{(i)} \), \( 1 \leq \alpha, \beta \leq m \), is an \( n \)-by-\( n \) Toeplitz matrix.

In [43], R. Chan, Nagy, and Plemmons considered preconditioners based on the level-1 and level-2 circulant approximations proposed by T. Chan and Olkin [59]. For \( j = 1, 2 \), they construct the level-\( j \) approximation \( c_j(A^{(i)}) \) to each \( A^{(i)} \); see §4.1. Then the preconditioner \( P_j \) is defined by

\[
P_j^* P_j = \sum_{i=1}^{k} c_j(A^{(i)})^* c_j(A^{(i)}).
\]

It can be shown that \( P_1 = (I \otimes F_n)^* \tilde{P}(I \otimes F_n) \), where \( \tilde{P} \) is a block upper triangular matrix with diagonal blocks. Thus, \( P_1 \) is block upper triangular with circulant blocks, and hence computing \( P_1^{-1} y \) involves the application of FFTs of size \( n \) together with backward solves involving \( \tilde{P} \). The computation of \( P_1^{-1} y \) is parallelizable and involves \( O(nm^2 + mn \log n) \) operations. For the level-2 preconditioner, we have \( P_2 = (F_m \otimes F_n)^* \Phi(F_m \otimes F_n) \), where \( \Phi \) is a diagonal matrix. Hence, \( P_2^{-1} y \) can be computed in \( O(mn \log(mn)) \) operations by using FFTs in both directions.

Convergence analysis has been given when all the blocks \( A^{(i)} \) in the matrix \( A \) are BTTB; i.e.,

\[
A^{(i)} = \begin{bmatrix}
A_0^{(i)} & A_1^{(i)} & \cdots & A_{2-m}^{(i)} & A_{1-m}^{(i)} \\
A_1^{(i)} & A_0^{(i)} & \cdots & A_{2-m}^{(i)} & \vdots \\
\vdots & \vdots & \ddots & \vdots & \vdots \\
A_{m-2}^{(i)} & \cdots & \cdots & A_0^{(i)} & A_{-1}^{(i)} \\
A_{m-1}^{(i)} & A_{m-2}^{(i)} & \cdots & A_1^{(i)} & A_0^{(i)}
\end{bmatrix}, \quad i = 1, \ldots, k,
\]

where each \( A_\mu^{(i)} \), \( |\mu| < m \), is a Toeplitz matrix with entries

\[
[A_\mu^{(i)}]_{\gamma,\delta} = a_{\mu,\gamma-\delta}^{(i)}, \quad 0 \leq \gamma, \delta < n.
\]
Notice that the generating function of $A^{(i)}$ is given by
\[
f^{(i)}(x, y) = \sum_{\mu=-\infty}^{\infty} \sum_{\nu=-\infty}^{\infty} a_{\mu, \nu}^{(i)} e^{-i\mu x} e^{-i\nu y} \quad \forall x, y \in [-\pi, \pi].
\]

**Theorem 4.6** (see R. Chan, Nagy, and Plemmons (1993) [43]). Let $A$ be given as in (4.14) with the entries of each $A^{(i)}$ satisfying
\[
\sum_{\mu=-\infty}^{\infty} \sum_{\nu=-\infty}^{\infty} |a_{\mu, \nu}^{(i)}| \leq K_i < \infty, \quad i = 1, \ldots, k.
\]
If one of the $f^{(i)}(x, y)$ is a positive function, then for any given $\epsilon > 0$ there exist positive integers $n_1$ and $n_2$ such that for all $m > n_1$ and $n > n_2$ at most $O(m) + O(n)$ eigenvalues of the matrix
\[
(P_j^* P_j)^{-1}(A^* A) - I, \quad j = 1, 2,
\]
have absolute values larger than $\epsilon$.

Numerical experiments in [43] show the effectiveness of the preconditioners proposed for block-Toeplitz least squares problems and actual image restoration problems. We remark that by using the new clustering result in Theorem 4.2 for BCCB preconditioners, we can conclude further that the preconditioned matrix $(P_2^* P_2)^{-1}(A^* A)$ will have a clustered spectrum around 1.

**4.4.2. Displacement-based preconditioners.** We first briefly review relevant definitions and results on displacement structure representation of Toeplitz matrices. We introduce the $n$-by-$n$ lower shift matrix $Z$, whose entries are zero everywhere except for 1's on the first subdiagonal. The displacement operator $\nabla$ is defined by
\[
\nabla A = A - ZAZ^*,
\]
where $\nabla A$ is called the displacement of $A$; see Chun and Kailath [63]. Let $L(w)$ denote the $n$-by-$n$ lower triangular Toeplitz matrix with the vector $w$ as its first column. Using these definitions, we have the following.

**Theorem 4.7** (see Chun, Kailath, and Lev-Ari (1987) [64]). An arbitrary $n$-by-$n$ matrix $A$ can be written in the form
\[
A = \sum_{i=1}^{\rho} L(u_i) L(v_i)^*,
\]
where $\rho = \text{rank}(\nabla A)$ and $u_i$ and $v_i$ are $n$-vectors.

The sum given in the theorem above is called the displacement representation of the given matrix $A$, and the scalar $\rho$ is called the displacement rank of $A$. Square Toeplitz matrices and Toeplitz-related matrices have small displacement rank [64]. For example, if $A$ is an $n$-by-$n$ Hermitian Toeplitz matrix, then
\[
A = L(x_+) L(x_+)^* - L(x_-) L(x_-)^*,
\]
where
\[
x_\pm = \left[ \frac{1}{2} (a_0 \pm 1), a_1, \ldots, a_{n-1} \right]^T.
\]
If $A$ is an $m$-by-$n$ ($m \geq n$) Toeplitz matrix, then $A^*A$ is in general not a Toeplitz matrix. However, $A^*A$ does have a small displacement rank $\rho \leq 4$ and a displacement representation in the form

$$
(4.15) \quad A^*A = \tilde{A} + L(y_1)L(y_1)^* - L(y_2)L(y_2)^*,
$$

where $\tilde{A}$ is a Hermitian Toeplitz matrix and

$$
y_1 = [0, a_{-1}, \ldots, a_{1-n}]^T \quad \text{and} \quad y_2 = [0, a_{m-1}, \ldots, a_{m-n+1}]^T.
$$

When the matrix $A$ is an $n$-by-$n$ Hermitian Toeplitz matrix, R. Chan, Nagy, and Plemmons [44] and Freund and Huckle [79] define the displacement preconditioner to be

$$
(4.16) \quad C = c(L(x_+))c(L(x_+))^* - c(L(x_-))c(L(x_-))^*.
$$

Clearly, $C$ is a Hermitian circulant matrix. R. Chan, Nagy, and Plemmons [44] proved that $C$ so defined in (4.16) is equal to the optimal circulant approximation of $A$; i.e., $C = c(A)$. In Freund and Huckle [79], displacement-based preconditioners for general matrices of low displacement ranks are also considered.

For general rectangular Toeplitz matrices, one can define, according to (4.15), the preconditioner

$$
c(\tilde{A}) + c(L(y_1))c(L(y_1))^* - c(L(y_2))c(L(y_2))^*.
$$

However, R. Chan, Nagy, and Plemmons [44] proved that the last term $L(y_2)L(y_2)^*$ in (4.15) is not significant as far as the conjugate gradient method is concerned. Thus, they defined the displacement preconditioner $C$ as

$$
C = c(\tilde{A}) + c(L(y_1))c(L(y_1))^*.
$$

Under the Wiener class assumptions on the generating function of $A$, clustering results have been established for $C$; see [44].

4.4.3. Strang’s generalized preconditioner. Recall that T. Chan’s circulant preconditioner is defined for general square matrices, not necessarily of Toeplitz form; see (2.6). Most circulant preconditioners including Strang’s do not enjoy this property. As a result, T. Chan’s preconditioners have been used extensively in solving non-Toeplitz systems in least squares problems.

Recently, R. Chan, Ng, and Plemmons [47] proposed a method to generalize the construction of Strang’s circulant preconditioner to arbitrary $n$-by-$n$ matrices $A$ and to employ this new circulant approximation in deconvolution applications in signal and image processing. For a general $n$-by-$n$ non-Toeplitz matrix $A = [a_{j,k}]$, we define its generalized Strang circulant preconditioner $S$ to be

$$
[S]_{j,\lfloor \frac{k}{2} \rfloor} = a_{j,\lfloor \frac{k}{2} \rfloor}, \quad 0 \leq j \leq n - 1;
$$

i.e., the $\lfloor \frac{k}{2} \rfloor$th column of $S$ is given by the $\lfloor \frac{n}{2} \rfloor$th column of $A$. Note that if $A$ is Toeplitz, then this definition coincides with that in (2.1).

This idea of constructing circulant preconditioners is similar to the forward–backward projection method used in [65] for image reconstruction. In [65], the authors tried to estimate the PSF involving a given, but not necessarily Toeplitz, matrix $A$ by forward projecting and
back projecting a one-pixel point source located at the center of the field of view. In matrix terms, the approximate PSF is just equal to the \( \left\lfloor \frac{n}{2} \right\rfloor \)th column of \( A \). The circulant matrix thus obtained was used in [65] as a preconditioner in the steepest descent method to speed up the convergence rate. The convergence rate of the preconditioned conjugate gradient method with the generalized Strang preconditioner was analyzed in [47], where again, under the Wiener class assumptions, the preconditioned matrix had a clustered spectrum.

**4.4.4. Inverse filter preconditioners.** Let us consider the convolution of a 1-d discrete signal \( x \) of length \( n \) with a convolution vector \( a \) of the form

\[
a = [a_{-m+1}, a_{-m+2}, \ldots, a_0, \ldots, a_{m-2}, a_{m-1}]^T.
\]

The resulting vector \( b \) is of length \( 2m + n - 2 \), and the convolution operation can be expressed in matrix notation as \( b = Ax \), where \( A \) is a column circulant matrix of the form

\[
A = 
\begin{bmatrix}
  a_{-m+1} & 0 & \cdots \\
  a_{-m+2} & a_{-m+1} & \ddots \\
  \vdots & \vdots & \ddots & 0 \\
  a_0 & \cdots & a_{-m+1} & \vdots \\
  \vdots & a_0 & \ddots & \vdots \\
  a_{m-2} & \cdots & \ddots & \vdots \\
  a_{m-1} & \cdots & a_0 & \vdots \\
  0 & a_{m-1} & \cdots & \vdots \\
  \vdots & \vdots & \ddots & \vdots \\
  0 & 0 & \cdots & a_{m-1}
\end{bmatrix}
\]

(4.17)

In deconvolution problems, the aim is to compute \( x \) given the matrix \( A \) and the vector \( b \). The first column of \( A \) is usually the discrete PSF of the blurring operator and is usually obtained empirically by blurring a point source image. In many applications of interest \( m \ll n \); see [116].

In signal and image processing, one common technique used to compute an approximate solution to the deconvolution problem is the inverse filter method. There the rectangular Toeplitz matrix \( A \) is embedded into an \( n \)-by-\( n \) circulant matrix \( C \), and the solution to the circulant system is used as an approximation to the solution of the Toeplitz system. This approach is attractive since \( n \)-by-\( n \) circulant systems can be solved in \( O(n \log n) \) operations by FFTs. However, the computed solution will not be very accurate.

Nagy, Plemmons, and Tørgerensen [145, 146] employed the inverse filter technique to construct a new preconditioner for Toeplitz least squares deconvolution problems. Their preconditioner \( M \) is based on the partition of the inverse of the circulant extension \( C \). Thus writing

\[
C^{-1} = F^* A F = \begin{bmatrix} M^* \\ M^* \\ \end{bmatrix},
\]

their preconditioner is defined to be \( M \), and they showed that \( M^* M \) is a good approximate inverse of \( A^* A \).
For 2-d deconvolution problems, one is still concerned with solving a least squares problem as in (4.13). But the matrix $A$ will be a block column circulant matrix with column circulant blocks, i.e.,

$$A = \begin{pmatrix}
A^{(-\ell+1)} & 0 \\
A^{(-\ell+2)} & A^{(-\ell+1)} \\
\vdots & \vdots & \ddots & 0 \\
A^{(0)} & \vdots & \cdots & A^{(-\ell+1)} \\
\vdots & A^{(0)} & \cdots & \vdots \\
A^{(\ell-2)} & \vdots & \cdots & \vdots \\
A^{(\ell-1)} & \vdots & \cdots & A^{(0)} \\
0 & A^{(\ell-1)} & \cdots & \vdots \\
\vdots & \vdots & \cdots & \vdots \\
0 & \cdots & \cdots & A^{(\ell-1)}
\end{pmatrix},$$

(4.18)

with each subblock $A^{(j)}$ being a $(2m+n-2)$-by-$n$ matrix of the form given by (4.17). We note that $A^* A$ will be an $n$-block-by-$n$ block-Toeplitz matrix with $n$-by-$n$ Toeplitz blocks. The construction of inverse filter preconditioners $M$ for the 2-d case proceeds as in the 1-d case. The matrix $A$ is extended to a BC matrix with circulant blocks. In [145, 146], Nagy, Plemmons, and Torgersen proved that the preconditioned matrices can also be written as a sum of the identity matrix and a matrix with rank dependent on $m$ and $\ell$.

4.4.5. Numerical example. The restoration of real images by using the preconditioned conjugate gradient algorithm with the preconditioners mentioned above has been carried out in [43, 44, 47, 145, 146]. Here we report the results found in [47]. We use the preconditioned conjugate gradient algorithm with the generalized Strang circulant preconditioner to remove the blurring in a 2-d image arising from atmospheric turbulence. The problem consists of a 256-by-256 image of an ocean reconnaissance satellite observed by a simulated ground-based imaging system together with a 256-by-256 image of a guide star observed under similar circumstances (see Figure 4.1). The data are provided by the Phillips Air Force Laboratory at Kirkland AFB, NM [25]. The imaging system estimates the atmospheric distortions using the natural guide star image. A wavefront sensor measures the optical distortions, which can then be digitized into a blurred image of the guide star pixel. To form the discrete PSF $a$, the rows of the blurred pixel image are stacked into a column vector. Then the PSF matrix $A$ is given in block form as in (4.18) with the stacked vector as its first column. We note that $A$ satisfies the Wiener class assumptions, since the guide star for the atmospheric imaging problem yields a Gaussian PSF [146].

In Figures 4.2 and 4.3, we present restorations without and with using the generalized Strang preconditioner described in §4.4.3. The regularization parameter $\mu$ in both cases is set to .018, as suggested in [145]. From the figures, we observe that when no preconditioner is used, an acceptable restoration is achieved after 50 iterations. Essentially, the same restoration is achieved in 10 iterations when preconditioning is used. Finally, we remark that roughly $0.61 \times 10^8$ floating point operations (done by Matlab) per iteration are used for the circulant-based deconvolution, while the count for no preconditioning is $0.50 \times 10^8$. 
Fig. 4.1. Observed image (left) and guide star image (right).

Fig. 4.2. Restored images without preconditioning: 10 (left) and 50 iterations (right).

Fig. 4.3. Restored image with preconditioning: 2 (left) and 10 iterations (right).
4.5. Applications to integral equations.

4.5.1. Inverse heat problem in $\mathbb{R}^m$. The inverse heat problem in $\mathbb{R}^m$ is the problem of recovering the initial data $\Upsilon(y)$ for all $y \in \mathbb{R}^m$ in

$$u(x,t) = \frac{1}{\sqrt{4\pi t}} \int_{\mathbb{R}^m} \exp\left(\frac{-(x-y)^2}{4t}\right) \Upsilon(y)dy$$

when for some $t > 0$, $u(x,t)$ is given for all $x \in \mathbb{R}^m$. Using discrete time and spatial sampling of the domain and sinc expansion for approximating the initial data, the problem is reduced to solving a linear system with block-Toeplitz coefficient matrices:

$$\begin{equation}
(A \otimes A \otimes \cdots \otimes A)\Upsilon = u,
\end{equation}$$

where $A$ is a Toeplitz matrix; see Gilliam, Martin, and Lund [83]. The generating function of the Toeplitz matrix $A$ is

$$f(\theta) = \exp\left(-\frac{\theta^2}{4\pi^2}\right), \quad \theta \in [-\pi, \pi],$$

which is positive and in the Wiener class. In [30], system (4.19) is preconditioned by $C \otimes C \otimes \cdots \otimes C$ with different circulant preconditioners $C$ discussed in §2.1. Numerical results in [30] show that for solving block-Toeplitz systems of order 65536, the time required by using the circulant preconditioned conjugate gradient method is half the time required by the direct Toeplitz solver available in the IMSL package.


$$\begin{equation}
y(t) + \int_{0}^{\infty} a(t-s)y(s)ds = g(t), \quad 0 \leq t < \infty,
\end{equation}$$

in which $a(t) \in L_1(\mathbb{R})$ and $g(t) \in L_2[0, \infty)$ are given functions, arise in a variety of practical applications in mathematics and engineering, such as linear prediction problems for stationary stochastic processes [90, pp. 145–146] and diffusion problems and scattering problems [90, pp. 186–189].

One way of solving (4.20) is by the projection method [85], where the solution $y(t)$ of (4.20) is approximated by the solution $y_\tau(t)$ of the finite-section equation

$$\begin{equation}
y_\tau(t) + \int_{0}^{\tau} a(t-s)y_\tau(s)ds = g(t), \quad 0 \leq t \leq \tau.
\end{equation}$$

It is shown in [85, Thm. 3.1] that

$$\lim_{\tau \to \infty} ||y_\tau - y||_{L_p[0,\tau]} = 0, \quad 1 \leq p < \infty.$$ 

The finite-section equation (4.21) can be solved numerically by either direct or iterative methods. For a fixed $\tau$, the finite-section operator $A_\tau$ defined by

$$\begin{equation}
(A_\tau x)(t) = \begin{cases} 
\int_{0}^{\tau} a(t-s)x(s)ds, & 0 \leq t \leq \tau, \\
0, & t > \tau 
\end{cases}
\end{equation}$$
is a compact operator. Therefore, the spectrum of the Wiener–Hopf operator \( \mathcal{I} + \mathcal{A}_\tau \) (where \( \mathcal{I} \) is the identity operator) is clustered around 1, and hence solving (4.21) by iterative methods such as the conjugate gradient method is less expensive than direct methods.

However, as \( \tau \) tends to infinity, the spectrum of this finite-section operator \( \mathcal{A}_\tau \) becomes dense in the spectrum of the half-line operator defined by (4.20), and hence the convergence rate of the conjugate gradient method deteriorates; see the numerical results in [86], for instance. A standard way of speeding up the convergence rate of the conjugate gradient method is to apply a preconditioner. Thus, instead of solving (4.21), one solves a preconditioned operator equation. We remark that there is a close relation between Wiener–Hopf integral equations and semi-infinite Toeplitz operators; see Gohberg and Fel’dman [85, p. 5]. The methodology of preconditioning “discrete” Toeplitz systems can be modified and applied to the “continuous” Wiener–Hopf equations.

In [86], Gohberg, Hanke, and Koltracht proposed using \textit{circulant integral operators} to precondition (4.21). Circulant integral operators are operators of the form

\[
(\mathcal{C}_\tau x)(t) = \int_0^\tau c_\tau(t-s)x(s)ds, \quad 0 \leq t \leq \tau,
\]

where \( c_\tau \) is a \( \tau \)-periodic conjugate symmetric function in \( L_1[-\tau, \tau] \); i.e.,

\[
c_\tau(t+\tau) = c_\tau(t) \quad \text{and} \quad c_\tau(-t) = \overline{c_\tau(t)} \quad \forall t \in [-\tau, \tau].
\]

We remark that \( \mathcal{C}_\tau \) is a compact, self-adjoint operator on \( L_2[-\tau, \tau] \). The preconditioned equation is given by

\[
(\mathcal{I} + \mathcal{C}_\tau)^{-1}(\mathcal{I} + \mathcal{A}_\tau)y_\tau(t) = (\mathcal{I} + \mathcal{C}_\tau)^{-1}g(t), \quad 0 \leq t \leq \tau.
\]

It has been shown in [86] that for large \( \tau \) the spectra of the circulant preconditioned operators \( (\mathcal{I} + \mathcal{C}_\tau)^{-1}(\mathcal{I} + \mathcal{A}_\tau) \) are clustered around 1.

**THEOREM 4.8** (see Gohberg, Hanke, and Koltracht (1994) [86]). If

\[
\lim_{\tau \to \infty} ||a - c_\tau||_{L_1[-\tau, \tau]} = 0,
\]

then for any given \( \epsilon > 0 \) there exist a positive integer \( \rho \) and a number \( \tau^* > 0 \) such that for all \( \tau > \tau^*, \mathcal{I} + \mathcal{A}_\tau \) are positive definite, and the spectra of \( (\mathcal{I} + \mathcal{C}_\tau)^{-1/2}(\mathcal{I} + \mathcal{A}_\tau)(\mathcal{I} + \mathcal{C}_\tau)^{-1/2} \) have at most \( \rho \) eigenvalues outside the interval \((1 - \epsilon, 1 + \epsilon)\).

Thus under assumption (4.25) the preconditioned conjugate gradient method converges superlinearly for large \( \tau \). In addition, Gohberg, Hanke, and Koltracht [86] proposed “wrap-round” and “optimal” circulant integral operators for Wiener–Hopf equations.

(1) **"Wrap-round" circulant integral operator** \( \mathcal{S}_\tau \). Given the operator \( \mathcal{A}_\tau \) as in (4.22), it is defined as

\[
(\mathcal{S}_\tau x)(t) = \begin{cases} 
\int_0^\tau s_\tau(t-s)x(s)ds, & 0 \leq t \leq \tau, \\
0, & t > \tau.
\end{cases}
\]

Here the function \( s_\tau \) is a \( \tau \)-periodic function defined by

\[
s_\tau(t) = a(t), \quad -\tau/2 \leq t \leq \tau/2.
\]
(2) "Optimal" circulant integral operator $F_\tau$. Given the operator $A_\tau$, it is defined as

$$(F_\tau x)(t) = \begin{cases} \int_0^t f_\tau(t-s)x(s)ds, & 0 \leq t \leq \tau, \\ 0, & t > \tau. \end{cases}$$

Here $f_\tau(t)$ is a $\tau$-periodic function defined by

$$f_\tau(t) = \left(\frac{\tau-t}{\tau}\right)a(t) + \left(\frac{t}{\tau}\right)a(t-\tau), \quad 0 \leq t \leq \tau.$$ 

Gohberg, Hanke, and Koltracht [86] showed that $F_\tau$ minimizes the Hilbert–Schmidt norm

$$|||A_\tau - C_\tau|||^2 = \int_0^\tau \int_0^\tau |a(t-s) - c_\tau(t-s)|^2 dsdt$$

over all circulant integral operators $C_\tau$.

We note that the above choices of $S_\tau$ and $F_\tau$ are continuous analogs of, respectively, Strang's and T. Chan's circulant preconditioners. In §2.2, a unified approach for constructing circulant matrix preconditioners for finite Toeplitz matrices is derived from the viewpoint of convolution products. Using the same viewpoint, R. Chan, Jin, and Ng [37] derived an easy and general scheme for constructing circulant integral preconditioners for Wiener–Hopf equations. For ease of presentation, let us denote by $\hat{q}$ the Fourier transform of any function $q$. The first step is to relate $c_\tau(t)$ in (4.23) to a sequence of conjugate symmetric functions $\{C_\tau(t)\}_\tau$, i.e.,

$$(4.27) \quad c_\tau(t) = \begin{cases} C_\tau(t)a(t) + C_\tau(t-\tau)a(t-\tau), & 0 \leq t \leq \tau, \\ C_\tau(t+\tau)a(t+\tau) + C_\tau(t)a(t), & -\tau \leq t \leq 0, \end{cases}$$

and translate the convergence requirement (4.25) on $c_\tau(t)$ to conditions on $\{C_\tau(t)\}_\tau$. Basically, $C_\tau(t)$ should be such that the convolution product $\hat{C}_\tau \ast \hat{a}$ converges uniformly to $\hat{a}$ on $\mathbb{R}$.

As an example, the "wrap-round" circulant integral operator $S_\tau$ can be constructed by setting $C_\tau$ in (4.27) to be the function

$$S_\tau(t) = \begin{cases} 1, & |t| \leq \tau/2, \\ 0, & |t| > \tau/2. \end{cases}$$

The operator $F_\tau$ can be also be derived by setting $C_\tau$ in (4.27) to be the function

$$F_\tau(t) = \begin{cases} \frac{\tau - |t|}{\tau}, & |t| \leq \tau, \\ 0, & |t| > \tau. \end{cases}$$

Moreover, R. Chan, Jin, and Ng proved that for sufficiently large $\tau$, if $\{C_\tau\}_\tau$ is uniformly bounded on the real line and the convolution product of $\hat{C}_\tau$ with $\hat{a}(\omega)$ converges to $\hat{a}(\omega)$ uniformly on $\mathbb{R}$, then the spectra of the circulant preconditioned operators are clustered around 1. They also showed that $\{C_\tau(t)\}_\tau$ can be derived easily from the Dirac delta function or from approximate convolution identities commonly used in Fourier analysis [190].

As in the Toeplitz matrix case, there are other ways of constructing operators as preconditioners for (4.21); see Ng and Lin [150], Ng, Lin, and R. Chan [151], and R. Chan and Lin.
In R. Chan and Lin [40], optimal and superoptimal circulant integral preconditioners are constructed for general integral equations of the second kind:

\[ y(t) + \int_0^\infty a(t, s)y(s)ds = g(t), \quad 0 \leq t < \infty. \]

Here \( a(t, s) \) is not necessarily a convolution kernel.

When (4.21) is discretized with the rectangular quadrature rule, we get an \( n \times n \) matrix system \( A\mathbf{y} = \mathbf{g} \), where \( A \) is a Toeplitz matrix. Here \( n \) is the number of quadrature points used in the discretization. We note that if the rectangular quadrature rule is used to discretize (4.24), then we get a matrix system

\[ (I + C)^{-1}(I + A)\mathbf{y} = (I + C)^{-1}\mathbf{g}, \tag{4.28} \]

where the matrices \((I + C)\) and \((I + A)\) are \( n \times n \) circulant and Toeplitz matrices, respectively. We see that (4.28) is basically a circulant-preconditioned Toeplitz system, which requires only \( O(n \log n) \) operations in each iteration.

The main drawback of using the rectangular rule is that the order of accuracy of the discretized solution \( \mathbf{y} \) depends only linearly on the number of quadrature points. Thus, in order to obtain a reasonably accurate solution for (4.21), a small step size must be used and hence the dimension of the resulting matrix system will be large. In order to obtain high order of accuracy, one can use higher-order quadrature rules such as the trapezoidal rule or Simpson’s rule, which have second and fourth orders of accuracy, respectively. In these cases, the discretization matrices of the Wiener–Hopf integral operators are non-Toeplitz matrices. Moreover, the corresponding discretization matrices of the circulant integral operators are in general not circulant, and therefore their inversion cannot be computed by using FFT. Hence, the cost per iteration of the preconditioned conjugate gradient method exceeds \( O(n \log n) \) operations.

Instead of constructing matrix preconditioners for these discretization matrices, we can consider preconditioners from the operator point of view. In [135], Lin, Ng, and R. Chan constructed preconditioners \( B_\tau \) for (4.21) such that the preconditioned operators \((\mathcal{I} - B_\tau)(\mathcal{I} + A_\tau)\) have clustered spectra and only \( O(n \log n) \) operations are required in each iteration of the preconditioned conjugate gradient method, even when higher-order quadrature rules are employed. Their idea is to use convolution operators rather than circulant operators to precondition \( \mathcal{I} + A_\tau \). More precisely, given \( A_\tau \), they constructed the convolution operator \( B_\tau \), whose kernel function \( b(t) \) satisfies

\[ \hat{b}(\omega) = \frac{\hat{a}(\omega)}{1 + \hat{a}(\omega)}. \tag{4.29} \]

Then they preconditioned (4.21) as

\[ (\mathcal{I} - B_\tau)(\mathcal{I} + A_\tau)\mathbf{y}_\tau(t) = (\mathcal{I} - B_\tau)\mathbf{g}(t), \quad 0 \leq t \leq \tau. \tag{4.30} \]

**Theorem 4.9** (see Lin, Ng, and R. Chan [135]). *Let \( a(t) \in L_1(\mathbb{R}) \) be conjugate symmetric and \( \hat{a}(\omega) \geq 0 \). Then for any given \( \epsilon > 0 \) there exist a positive integer \( \rho \) and a real number \( \tau^* > 0 \) such that for all \( \tau > \tau^* \) the spectrum of \((\mathcal{I} - B_\tau)(\mathcal{I} + A_\tau)\) has at most \( \rho \) eigenvalues outside the interval \((1 - \epsilon, 1 + \epsilon)\).*

According to this theorem, if we apply the conjugate gradient method to solve (4.30), the convergence rate will be superlinear for large \( \tau \).
In the following, we test the effectiveness of different integral operator preconditioners by using the kernel function $a(t) = (1 + t^2)^{-1}/\mu$. In practical applications, the parameter $\mu$ is the regularization parameter and is usually a small positive number. In the test, we arbitrarily set $\mu = 0.01$. To discretize (4.29) we partition the interval $[-\tau, \tau]$ into $2n$ equal subintervals of step size $h$ and then compute $\hat{a}(\omega)$ and $\hat{b}(\omega)$ by using the formula

$$\hat{a}(\omega) \approx h \sum_{k=-n}^{n} a(kh)e^{-i\omega kh},$$

which can be computed easily by FFTs. The right-hand side function $g(t)$ is chosen such that the corresponding solution for the Wiener–Hopf equation (4.20) is

$$y(t) = \begin{cases} (16 - t)^2, & 0 \leq t \leq 16, \\ 0, & t > 16. \end{cases}$$

(4.31)

We use the same random vector as an initial guess for all preconditioners, and the tolerance is set to $10^{-6}$. All computations are done by Matlab. We define the error of the numerical solution:

$$e \equiv \left\{ h \sum_{j=0}^{n} |\tilde{y}(jh) - y(jh)|^2 \right\}^{1/2} \approx \left\{ \int_0^\tau |\tilde{y}(t) - y(t)|^2 dt \right\}^{1/2},$$

where $\{\tilde{y}(jh)\}_{j=0}^{n}$ is the computed solution and $y(t)$ is the true solution given by (4.31).

Table 4.5 gives the number of megaflops used to achieve a given accuracy $\epsilon$ (i.e., the error of the numerical solution $e \leq \epsilon$) for different quadrature rules and preconditioners. The symbol $B$ denotes that we are solving (4.30) with the quadrature rule listed. The symbol $S$ denotes that we are solving (4.24) with the “wrap-round” circulant integral operator $S_{r}$ defined in (4.26), and that $S_{r}$ is discretized by the rectangular rule (so as to make it a circulant matrix), but the operator $A_{r}$ is still discretized according to the rule listed in the table. The discretization matrix for $A_{r}$ will be a product of a Toeplitz matrix and a diagonal matrix, where the diagonal matrix depends on the discretization rule used. The symbol $I$ denotes that (4.21) is solved without using any preconditioner and discretized according to the rule listed in the table. ** denotes that the corresponding number exceeds 1000.

From the table, we see that without using any preconditioner the number of flops for achieving the required accuracy is very large. We note that for the trapezoidal rule the Toeplitz structure is disturbed only in the two rows corresponding to the boundary. Therefore, the circulant preconditioner works well for the trapezoidal case. However, it does not work well...
if \( \mathcal{A} \) is discretized using Simpson's rule. The performance of our proposed preconditioner for Simpson’s rule is the best one in terms of the accuracy and the computational work. The accuracy of the computed solution depends only on the quadrature rule used in discretizing \( \mathcal{A} \). However, the convergence rate of the preconditioned systems and the cost per iteration of the PCG method depend on how we discretize the preconditioning operators. Therefore, it is advantageous to use a higher-order quadrature rule to discretize the operator equation because of the accuracy concern. But to speed up the convergence rate of the method and minimize the cost per iteration, one needs to use our proposed preconditioner rather than circulant ones.

4.6. Applications to time series analysis.

4.6.1. Finite impulse response filters. Least squares estimations have been used extensively in a wide variety of scientific applications, such as equalization [84, p. 139], system identification [140, 141, 142], adaptive signal processing [3, 170, 99], active noise control [59], and speech processing [84, p. 343]. In these applications, we usually need to estimate the transmitted signal from a sequence of received signal samples or model an unknown system by using a linear system model.

To present the problem properly, let us introduce some terminology used in signal processing. Let \( x(i) \) be a discrete-time stationary zero-mean complex-valued process (see Fuller [82, pp. 10–11]). A finite impulse response (FIR) linear filter of order \( n \) is of the form

\[
\tilde{d}(i) = \sum_{k=1}^{n} b_k x(i - k + 1),
\]

where \( \tilde{d}(i) \) is the filter output based on the data \( \{x(k)\}_{k=i}^{i+n-1} \) and \( \{b_k\}_{k=1}^{n} \) are the impulse responses of the filter. The difference between the desired response \( d(i) \) of the process and the filter output \( \tilde{d}(i) \) is called the estimation error of order \( n \). Since we are interested in estimating the desired response based on the input measurements, the impulse responses \( \{b_k\}_{k=1}^{n} \) should be chosen to make the estimation error as small as possible.

For the case of known statistics, i.e., the autocovariances of the stationary process are known, the optimal least mean squares predictor coefficients \( \{b_k\}_{k=1}^{n} \) are given by the solution of the linear system of equations

\[
A b = d;
\]

see Giordano and Hsu [84, pp. 41–43]. Here \( A \) is an \( n \)-by-\( n \) Hermitian Toeplitz matrix and its entries \( \{a_j\}_{j=0}^{n-1} \) are the autocovariances of a discrete-time stationary process and are given by

\[
a_k = \mathcal{E}[x(j)x(j-k)],
\]

where \( \mathcal{E} \) is the expectation operator.

The matrix \( A \) is called the covariance matrix of the stationary process, and the Toeplitz system (4.32) is commonly called the Yule–Walker equations; see Yule [197].

For a discrete-time stationary process, if the autocovariances of the process are absolutely summable, i.e., \( \sum_{k=-\infty}^{\infty} |a_k| < \infty \), then the function \( f(\theta) \) with \( a_k \) as Fourier coefficients is called the spectral density function of the stationary process; see [21, p. 118]. The covariance matrix \( A \) is then a Toeplitz matrix generated by \( f(\theta) \). As examples, we consider the following stationary processes.

1. **First-order autoregressive process AR(1)** [160, p. 238]: The process is given by

\[
x(t) = \rho x(t - 1) + v(t),
\]
where \( \{v(t)\} \) is a white noise process with variance \( \eta^2 \). The autocovariances of the process are given by

\[
 a_k = \frac{\eta^2 \rho^{(|k|)}}{1 - \rho^2}, \quad k = 0, \pm 1, \pm 2, \ldots,
\]

where \( |\rho| < 1 \). The corresponding spectral density function \( f(\theta) \) is given by

\[
f(\theta) = \frac{\eta^2}{2\pi(1 - 2\rho \cos \theta + \rho^2)}.
\]

The covariance matrix is a scalar multiple of the Kac–Murdock–Szegö matrix; see [122].

2. First-order moving average process MA(1) [21, p. 123]: The process is given by

\[
x(t) = v(t) + \chi v(t - 1),
\]

where \( |\chi| < 1 \) and \( \{v(t)\} \) is a white noise process with variance \( \eta^2 \). The autocovariances of the process are given by

\[
 a_k = \begin{cases} 
 \eta^2(1 + \chi^2), & k = 0, \\
 \eta^2 \chi, & k = 1, \\
 0, & \text{otherwise.}
\end{cases}
\]

We see that the covariance matrix is a tridiagonal Toeplitz matrix with

\[
f(\theta) = \frac{\eta^2}{2\pi}(1 + 2\chi \cos \theta + \chi^2).
\]

If one assumes that the spectral density function of the stationary process exists and is positive (that can be guaranteed by the causality of the process [21, p. 85]), then the Yule–Walker equations (4.32) can be solved in \( O(n \log n) \) operations by using the preconditioned conjugate gradient method with circulant preconditioners discussed in §2.1; see [149]. We remark that all the above results are derived deterministically. In the least squares estimation algorithms, we always deal with data samples from random processes and the convergence rate should be considered in a probabilistic way.

4.6.2. Least squares filters. We note that in practical cases no prior knowledge is usually available on the autocovariances of the process. If \( m \) data samples have been taken, then all the information we have is contained in the finite number of data points \( \{x(t)\}_{t=1}^m \). In this case, one can still formulate a well-defined least squares prediction problem by estimating the autocovariances from the data samples \( \{x(t)\}_{t=1}^m \) with various types of windowing methods, such as the correlation, covariance, prewindowed, and postwindowed methods; see, for instance, Giordano and Hsu [84, pp. 65–66]. The least squares estimators can then be found by solving for the \( n \)-vector \( b \) in

\[
\min ||d - Ab||_2.
\]

Here \( A \) is an \( m \)-by-\( n \) Toeplitz matrix with full column rank \( n \), obtained by applying various types of windowing methods on the data samples \( \{x(t)\}_{t=1}^m \). The solution \( b \) of (4.34) can be obtained by solving the normal equation \( (A^*A)b = A^*d \). We note that if the correlation
method is employed, the normal matrix $A^* A$ is Toeplitz. The other three windowing methods lead to non-Toeplitz normal matrices $A^* A$ of the form

\begin{equation}
A^* A = \tilde{A} - L^* L - U^* U,
\end{equation}

where $\tilde{A}$ is an $n$-by-$n$ Toeplitz matrix and $L$ and $U$ are lower triangular and upper triangular Toeplitz matrices, respectively (cf. (4.15)).

To prove convergence, the following practical signal processing assumptions of the random process were made in [149, 148].

- The process is stationary with constant mean $\mu$.
- The spectral density function of the process is positive and in the Wiener class.
- There exist positive constants $\beta_1$ and $\beta_2$ such that

$$\Var \left( \frac{1}{m} \sum_{j=1}^{m-k} x(j) \right) \leq \frac{\beta_1}{m}, \quad k = 0, 1, 2, \ldots, m - 1$$

and

$$\Var \left( \frac{1}{m} \sum_{j=1}^{m-k} [x(j) - \mu][x(j+k) - \mu] \right) \leq \frac{\beta_2}{m}, \quad k = 0, 1, 2, \ldots, m - 1.$$  

We note that the positiveness of the spectral density function can be guaranteed by the causality of the process [21, p. 85], whereas the absolute summability of the autocovariances can be assured by the invertibility of the process [21, p. 86]. With these assumptions, we proved that the spectra of the preconditioned matrices $c(\tilde{A})^{-1}(A^* A)$ are clustered around 1 with probability 1, provided that a sufficiently large number of data samples are taken.

**Theorem 4.10** (see Ng and R. Chan (1994) [149]). *Let the discrete-time process satisfy the above assumptions. Then for any given $\epsilon > 0$ and $0 < \eta < 1$ there exist positive integers $\rho_1$ and $\rho_2$ such that for $n > \rho_1$ the probability that at most $\rho_2$ eigenvalues of the matrix $I - c(\tilde{A})^{-1}(A^* A)$ have absolute value greater than $\epsilon$ is greater than $1 - \eta$, provided that $m = O(n^{3+\nu})$ with $\nu > 0$.

Hence, when we apply the conjugate gradient method to the preconditioned system, the method converges superlinearly with probability 1. Since the data matrix $A$ is an $m$-by-$n$ rectangular Toeplitz matrix, the normal equations and the circulant preconditioner can be formed in $O(m \log n)$ operations. Once they are formed, the cost per iteration of the preconditioned conjugate gradient method is $O(n \log n)$ operations. Therefore, the total work of obtaining the predictor coefficients to a given accuracy is of $O((m+n) \log n)$.

**4.6.3. Linear-phase least squares filters.** Besides FIR linear filters, FIR linear-phase filters are also commonly used in signal processing. Such filters are especially important for applications where frequency dispersion due to nonlinear phase is harmful, such as speech processing. In this case, the impulse responses can be found by solving the Toeplitz-plus-Hankel least squares problem

$$\min \|d - (A + H)b\|_2,$$

where $A + H$ is a rectangular Toeplitz-plus-Hankel matrix; see [140, 141, 142, 194, 109, 147]. By exploiting the structure of the normal equations matrices, it can be written as

$$(A + H)^*(A + H) = \tilde{A} + \tilde{H} - V_1 - V_2 - V_3 - V_4,$$
where $\tilde{A}$ is a Toeplitz matrix, $\tilde{H}$ is a Hankel matrix, and $\{V_i\}_{i=1}^4$ are non-Toeplitz and non-Hankel matrices. In [147], $c(\tilde{A})$ is used as a preconditioner for the problem. Under the same assumptions as in Theorem 4.10, it has been shown that the spectra of the matrices $\tilde{H}$ and $\{V_i\}_{i=1}^4$ are all clustered around 0. Hence, the spectrum of the preconditioned matrix $c(\tilde{A})^{-1}(A + H)^*(A + H)$ will also be clustered around 1 with probability 1.

Ku and Kuo [133] have also proposed a preconditioner for Toeplitz-plus-Hankel systems. Their approach is basically to take circulant approximations of the Toeplitz matrix and the Hankel matrix separately and then combine them to form a preconditioner. The motivation behind Ng's preconditioner is that the Toeplitz matrix $\tilde{A}$ is the sample autocorrelation matrix, which intuitively should be a good estimation to the autocorrelation matrix of the discrete-time stationary process, provided that a sufficiently large number of data samples are taken. Moreover, under practical signal processing assumptions, the spectrum of the Hankel matrix $\tilde{H}$ is clustered around 0. Hence, it suffices to approximate the Toeplitz part by circulant preconditioners.

4.6.4. Recursive least squares filters. So far we have discussed only block-processing-type least squares estimations; i.e., data samples are collected over a finite time interval and the resulting linear system is solved. Recently, Plemons [159] and Ng and Plemons [152, 153] proposed using circulant preconditioners for the recursive least squares (RLS) estimations, which are used extensively in many signal processing and control applications. They considered the RLS computations where the data matrices are assumed to have a Toeplitz (displacement) structure. Their new algorithm computes least squares estimators recursively by using sliding data windows involving multiple (rank $k$) updating and downdating computations for superior tracking capabilities. When $A(t)$ is an $m$-by-$n$ rectangular data matrix ($m$ is the length of the sliding window) with full column rank, then the least squares estimator $b(t)$ at step $t$ can be obtained by solving the normal equations

$$A(t)^*A(t)b(t) = A(t)^*d(t).$$

We note that although $A(t)^*A(t)$ is generally not Toeplitz, it can still be written in the form

$$A(t)^*A(t) = \tilde{A}(t) - L(t)^*L(t) - U(t)^*U(t),$$

where $\tilde{A}(t)$ is Toeplitz and $L(t)$ and $U(t)$ are lower triangular and upper triangular Toeplitz matrices (cf. (4.15) and (4.35)). In [152, 153], Ng and Plemons employed the preconditioned conjugate gradient method with circulant preconditioners to solve such systems at each step. In the case of point-processing ($k = 1$), the method requires $O(n \log n)$ operations per adaptive filter input, where $n$ is the number of least squares estimators. In the case of block-processing ($k \geq n$), the method requires only $O(\log n)$ operations per adaptive filter input. These FFT-based iterative RLS algorithms, with reasonable complexity for computing least squares estimators recursively, may also avoid some of the instability problems associated with direct fast RLS methods.

In the following, we test the convergence performance of the FFT-based sliding window RLS algorithm and compare it with fast transversal filter algorithms [99, pp. 586–600] and standard RLS [99, p. 485]. We remark that they are $O(n)$ and $O(n^2)$ algorithms, respectively. An exponential weighting factor $\gamma$ is generally used in these two algorithms. The inverse of $1 - \gamma$ is approximately a “measure” of the memory of the algorithm. Therefore, the length of the sliding window $m$ used in the FFT-based sliding window RLS algorithm is related to $\gamma$ by $m \approx 1/(1 - \gamma)$. 


Figure 4.4 is a block diagram for the implementation of the FFT-based RLS algorithm in an adaptive FIR system identification model. The input signal $x(t)$ drives the unknown system to produce the output sequence $d(t)$. We model the unknown system as an FIR filter. As an input stochastic process, we used the first-order autoregressive AR(1) process given by

$$x(t) + \rho x(t - 1) = v(t),$$

where $\rho = -0.9999$ and $\{v(t)\}$ is a white noise process with variance $\eta^2 = 1$. We note that the spectral density function of the process is positive and in the Wiener class; see (4.33). The reference (unknown) system is an $n$th-order linear-phase FIR filter with uncorrelated Gaussian white noise added. The FIR’s impulse responses $\{b_k\}_{k=1}^n$ of the unknown system used are

$$b_k = 1.1 - \frac{|2k - n - 1|}{n - 1}, \quad k = 1, 2, \ldots, n.$$  

In the numerical tests, the stopping criterion for the preconditioned conjugate gradient method is when the $\ell_2$ norm of the residual vector is less than $10^{-7}$. All the computations were done by Matlab.

Figures 4.5, 4.6, and 4.7 show the prior average least squares error and the $\ell_2$ norm of the weight error vector for different algorithms when different levels of variances of Gaussian white noise are added to the reference system. The size $n$ of the filters used is 32. We see from the figures that the fast transversal filter algorithm does not converge. However, both the standard RLS and the FFT-based sliding window algorithms converge very fast.

5. Concluding remarks. In 1986, Strang addressed the question of whether iterative methods can compete with direct methods for solving symmetric positive-definite Toeplitz systems. The answer has turned out to be an unqualified yes. The conjugate gradient method coupled with a suitable preconditioner can solve a large class of $n$-by-$n$ Toeplitz systems in $O(n \log n)$ operations, as compared to the $O(n \log^2 n)$ operations required by fast direct Toeplitz solvers. This paper summarizes some of the developments of this iterative method over the past few years. Applications of the method to partial differential equations, queueing networks, integral equations, image restoration, and time series analysis are also given. The results show that the method in some instances works better than traditional methods used specifically for these problems.
A possible direction of future work is the application of the method to other areas where solutions of Toeplitz or Toeplitz-related systems are sought, such as control theory and image processing. For instance, the restoration of images in nonlinear space-invariant systems involves the solutions of Toeplitz-like systems (see [178]), and the total variation method for restoring noisy, blurred images involves operators that are the sum of second-order elliptic
operators and Toeplitz operators [163, 189]. Also, in the numerical solution of the biharmonic equation by conformal mapping we must solve Toeplitz-plus-Hankel systems [34].

Other directions consist of a thorough comparison of the different preconditioners proposed and, more importantly, with fast direct Toeplitz solvers (cf. [113, 136]), and the use of other iterative methods, such as multigrid methods, for Toeplitz systems; see [74, 75].

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