# Quadrature Methods for Periodic Singular and Weakly Singular Fredholm Integral Equations 

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#### Abstract

High-accuracy numerical quadrature methods for integrals of singular periodic functions are proposed. These methods are based on the appropriate Euler-Maclaurin expansions of trapezoidal rule approximations and their extrapolations. They are subsequently used to obtain accurate quadrature methods for the solution of singular and weakly singular Fredholm integral equations. Throughout the development the periodic nature of the problem plays a crucial role. Such periodic equations are used in the solution of planar elliptic boundary value problems such as those that arise in elasticity, potential theory, conformal mapping, free surface flows, etc. The use of the quadrature methods is demonstrated with numerical examples.


KEY WORDS: Fredholm integral equations; singular integral equations; quadrature methods; boundary integrals.

## 1. INTRODUCTION

In this work we shall present Romberg-type integration formulas for the numerical evaluation of singular integrals whose integrands are periodic. We shall subsequently use these integration formulas to develop new highorder quadrature methods for the numerical solution of singular and weakly singular Fredholm integral equations of the first and second kinds

$$
\begin{equation*}
\omega f(t)+\int_{a}^{b} K(t, x) f(x) d x=g(t), \quad a \leqslant t \leqslant b \tag{1.1}
\end{equation*}
$$

[^0]where $\omega=0$ and $\omega=1$ for first and second kinds, respectively, and the corresponding eigenvalue problems
\[

$$
\begin{equation*}
f(t)+\lambda \int_{a}^{b} K(t, x) f(x) d x=0, \quad a \leqslant t \leqslant b \tag{1.2}
\end{equation*}
$$

\]

subject to the following conditions:
(1) The kernel $K(t, x)$ in (1.1) and (1.2) is periodic both in $t$ and in $x$ with period $T=b-a$, and for $t \neq x$ is differentiable in $t$ and $x$ as many times as needed. We furthermore assume that it is of the form

$$
\begin{aligned}
K(t, x)= & \sum_{k=1}^{N} W_{k}(t, x)\left|\phi_{k}(t)-\phi_{k}(x)\right|^{\beta_{k}}\left[\log \left|\widetilde{\phi}_{k}(t)-\widetilde{\phi}_{k}(x)\right|\right]^{q_{k}} \\
& +\frac{\hat{W}_{1}(t, x)}{\hat{\phi}_{1}(t)-\hat{\phi}_{1}(x)}+\hat{W}_{2}(t, x)
\end{aligned}
$$

where $\beta_{k}$ are real numbers satisfying $\beta_{k}>-1$ and $q_{k}$ are nonnegative integers. The functions $W_{k}(t, x)$ and $\hat{W}_{j}(t, x)$ are periodic in $t$ and $x$ with period $T$, and differentiable in $t$ and $x$ as many times as needed. Similarly, the functions $\phi_{k}(t), \widetilde{\phi}_{k}(t)$, and $\hat{\phi}_{1}(t)$ are periodic in $t$ with period $T$, and are differentiable as many times as needed. It is assumed that $W_{k}(t, t) \neq 0$ whenever $W_{k}(t, x) \not \equiv 0$ for some $k$. Similarly, it is assumed that $\hat{W}_{1}(t, t) \neq 0$ whenever $\hat{W}_{1}(t, x) \not \equiv 0$, and in this case the integrals in (1.1) are to be taken as Cauchy principal value integrals. When $W_{k} \equiv 0$ for all $k$ and $\hat{W}_{1} \not \equiv 0$ the integral equations in (1.1) and (1.2) are called singular, and when $W_{k} \not \equiv 0$ for some $k$ and $\hat{W}_{\mathrm{I}} \equiv 0$, they are called weakly singular. It is worth mentioning that the more important cases of $K(t, x)$ that arise in applications are

$$
\begin{aligned}
& K(t, x)=W_{1}(t, x) \log \left|\widetilde{\phi}_{1}(t)-\tilde{\phi}_{1}(x)\right|+\hat{W}_{2}(t, x) \quad \text { (weakly singular) } \\
& K(t, x)=\frac{\hat{W}_{1}(t, x)}{\hat{\phi}_{1}(t)-\hat{\phi}_{1}(x)}+\hat{W}_{2}(t, x) \quad \text { (singular) }
\end{aligned}
$$

and combinations of the two.
(2) The right-hand side, $g(t)$, in (1.1) is periodic in $t$ with period $T=b-a$, and is assumed to be differentiable as many times as needed.
(3) We assume that the solution $f(t)$ exists uniquely, is periodic in $t$ with period $T=b-a$, and is differentiable as many times as needed. [A heuristic argument to justify the assumption about the smoothness of $f(t)$ will be given later in this section.]

We note that integral equation formulations of two-dimensional boun-
dary value problems over domains with rectifiable boundaries naturally result in Fredholm integral equations (single or coupled) with singular or weakly singular (in general logarithmic) and periodic kernels and periodic input functions and solutions. This class of problems is by no means limited. Laplace's equation and the biharmonic equation, problems in elasticity, conformal mapping, free surface flows, etc., are but a few examples. For example, the conformal mapping problem can be solved via the equation of Symm (1966), which is a Fredholm integral equation of the first kind with a logarithmic kernel (see Example 3 in this work for additional references). The same problem can also be solved via Theodorsen's equation (see Gaier, 1964, p. 65), which is a nonlinear singular Fredholm equation with the Hilbert kernel as its kernel. Similarly, the solution of the biharmonic equation can be obtained with the help of the solution to a weakly singular coupled system of two Fredholm equations (see, for example, Jaswon and Symm, 1977, Chaps. 9 and 15). In all cases, the kernels $K(t, x)$ in the resulting integral equations are singular for $x=t$, but smooth otherwise, provided the boundary of the domain over which the original problems are defined is smooth too. The solutions seem to be smooth in general, provided the input functions are smooth in addition. For some cases of interest this can be shown rigorously.

Integral equation formulations of practical problems such as those mentioned in the previous paragraph are used very widely in many branches of engineering and fall in the category of "boundary integral equation methods."

One of the methods for solving (1.1) and (1.2) numerically is the quadrature method (see Baker, 1977, Chap. 4, Section 3), in which one replaces the integral $\int_{a}^{b} K(t, x) f(x) d x$ by a numerical quadrature formula, whose abscissas are $x_{j}, j=1, \ldots, n$, with $t=x_{i}, i=1, \ldots, n$, then replaces the $f\left(x_{j}\right)$ by their approximation $\bar{f}_{j}$, and finally solves the resulting system of linear equation for the $\tilde{f}_{j}$. Obviously, the accuracy of this method depends on the accuracy of the numerical quadrature formula being used, which in turn depends on the analytic properties of both the kernel $K(t, x)$ and the solution $f(t)$ over $[a, b]$. It can be said, in general, that whenever $K(t, x)$ is weakly singular or singular, the solution $f(t)$ will be singular at the end points $a$ and $b$. The singularity structure of $f(t)$ may be complicated and difficult to determine; see MacCamy (1958) and Graham (1982) for some general results on this problem. When $K(t, x), g(t)$, and $f(t)$ are (periodic) as assumed in the present work, then $a$ and $b$ in (1.1) and (1.2) can be replaced by $a^{\prime}$ and $b^{\prime}$, respectively, where $b^{\prime}-a^{\prime}=T$. If we now assume that $f(t)$ has singularities at $a$ and $b$, then it should be singular at $a^{\prime}$ and $b^{\prime}$ and hence at all $t$. As a result we conclude, heuristically, that $f(t)$ cannot have any singularities, and this is the assumption that we have made above.

In relation to the kernel $K(t, x)$, it is important to note that it can also be reexpressed in the form

$$
\begin{equation*}
K(t, x)=\sum_{k=1}^{M} H_{k}(t, x)|t-x|^{\alpha_{k}}(\log |t-x|)^{p_{k}}+\frac{\hat{H}_{1}(t, x)}{t-x}+\hat{H}_{2}(t, x) \tag{1.3}
\end{equation*}
$$

This is the form to be used in the developments throughout the remainder of this work. Here $\alpha_{k}$ are real numbers satisfying $\alpha_{k}>-1$ and $p_{k}$ are nonnegative integers. The $H_{k}(t, x)$ and $\hat{H}_{j}(t, x)$ have the same properties as the $W_{k}(t, x)$ and $\hat{W}_{j}(t, x)$, respectively, with the exception that they are not required to be periodic, neither in $t$ nor in $x$. This is so as none of the terms $|t-x|^{\alpha_{k}}(\log |t-x|)^{p_{k}}$ and $(t-x)^{-1}$ is. It is only the combination given on the right-hand side of (1.3) that is required to be periodic. Furthermore, as will become clear later in this work, full knowledge of the $H_{k}(t, x)$ and $\hat{H}_{j}(t, x)$ is not needed. What is needed, in general, is $H_{k}(t, t), 1 \leqslant k \leqslant M$, and $\hat{H}_{2}(t, t)$. For the particular case of the singular integral equations, however, neither $\hat{H}_{1}(t, t)$ nor $\hat{H}_{2}(t, t)$ is needed.

Finally, the task of determining the appropriate $H_{k}(t, t)$ and $\hat{H}_{2}(t, t)$ can be accomplished by simply expanding $K(t, x)$ for $x \rightarrow t$. As an example, consider $K(t, x)=\log |z(t)-z(x)|$, where $z(\tau)$ is a complex periodic function of $\tau$ with period $T$ (see Example 3 in this work). We can express $K(t, x)$ as in (1.3) as follows:

$$
K(t, x)=\log |t-x|+\log \left|\frac{z(t)-z(x)}{t-x}\right|
$$

i.e.,

$$
H_{1}(t, x)=1 \quad \text { and } \quad \hat{H}_{2}(t, x)=\log \left|\frac{z(t)-z(x)}{t-x}\right|
$$

Thus, $H_{1}(t, t)=1$. Expanding $z(x)$ about $x=t$, and letting $x \rightarrow t$, we also obtain $\hat{H}_{2}(t, t)=\log \left|z^{\prime}(t)\right|$. Note that $\hat{H}_{2}(t, x)$ is not periodic, although $K(t, x)$ is.

We now give a brief outline of the developments of the present work. Let

$$
\begin{equation*}
x_{j}=a+j h, \quad h=(b-a) / n, \quad n \text { a positive integer } \tag{1.4}
\end{equation*}
$$

Using the Euler-Maclaurin expansion for smooth integrands, and their extension to integrands having end-point singularities (see Navot, 1961, 1962), in the next section we derive Euler-Maclaurin expansions essentially for the integrals $\int_{a}^{b} K(t, x) f(x) d x$, with $K(t, x)$ and $f(x)$ as described
previously. We actually derive asymptotic expansions, for $h \rightarrow 0$, for the differences

$$
\begin{equation*}
\Delta(t, h)=I[t ; f]-I_{n}[t ; f] \tag{1,5}
\end{equation*}
$$

where

$$
\begin{equation*}
I[t ; f]=\int_{a}^{b} K(t, x) f(x) d x \tag{1.6}
\end{equation*}
$$

and $I_{n}[t ; f]$, the approximation to $I[t ; f]$, has the form

$$
\begin{equation*}
I_{n}[t ; f]=\sum_{j=1}^{n} w_{n}\left(t, x_{j}\right) f\left(x_{j}\right) \tag{1.7}
\end{equation*}
$$

with the $x_{j}$ as defined in (1.4), such that $t$ is one of the points $x_{j}$, and is being held fixed, and $w_{n}\left(t, x_{j}\right)=h K\left(t, x_{j}\right)$ for $x_{j} \neq t$, and $w_{n}(t, t)$ depend on the type of singularity that $K(t, x)$ has for $x=t$. For the case in which $K(t, x)=H_{1}(t, x) \log |t-x|+\hat{H}_{2}(t, x)$ we have

$$
\begin{aligned}
& w_{n}(t, t)=h\left[\hat{H}_{2}(t, t)+\log \left(\frac{h}{2 \pi}\right) H_{1}(t, t)\right] \quad \text { and } \\
& \Delta(t, h) \sim \sum_{i=1}^{\infty} \beta_{i} h^{2 i+1} \\
& \text { as } h \rightarrow 0
\end{aligned}
$$

where the $\beta_{i}$ are independent of $h$ and depend only on $t$. Consequently, the error in the quadrature formula $I_{n}[t ; f]$ is $O\left(h^{3}\right)$ as $h \rightarrow 0$. Similarly, for the case in which $K(t, x)=H_{1}(t, x)|t-x|^{\alpha}+\hat{H}_{2}(t, x)$, we have

$$
\begin{aligned}
& w_{n}(t, t)=h\left[\hat{H}_{2}(t, t)-2 \zeta(-\alpha) H_{1}(t, t) h^{\alpha}\right] \text { and } \\
& \Delta(t, h) \sim \sum_{i=1}^{\infty} \beta_{i} h^{2 i+\alpha+1} \quad \text { as } h \rightarrow 0,
\end{aligned}
$$

where, again, the $\beta_{i}$ are independent of $h$ and depend only on $t$. As a result, the error in $I_{n}[t ; f]$ is $O\left(h^{3+\alpha}\right)$ as $h \rightarrow 0$. Here $\zeta(s)$ is the Rieman zeta function.

Using the asymptotic expansions for $\Delta(t, h)$, in Section 3 we derive Romberg-type numerical quadrature formulas for $I[t ; f]$, thus increasing the accuracy of $I_{n}[t ; f]$ by as many orders of magnitude as we wish. In Section 4 we devise quadrature methods for the integral equations mentioned above that are based on these Romberg-type formulas. In Section 5 we illustrate the efficiency of our new quadrature methods with numerical examples. In Section 6 we review some quadrature methods that have been proposed in the past and bear some relation to the ones derived in the present work.

Note. The treatments of singular and weakly singular integral equations are not separated from each other throughout this paper. The reader interested in the treatment of singular integral equations could follow it easily by going through Theorems $2.1,2.4,2.7 \mathrm{a}, 3.1,3.2$, and the first part of Section 4, without having to consider the rest of this work.

## 2. EULER-MACLAURIN EXPANSIONS

The notation described below will be used throughout the remainder of this work.

Let $x_{j}=a+j h, j=0,1, \ldots, n, h=(b-a) / n$, where $n$ is a positive integer. Let $t \in(a, b)$ be fixed and $t \in\left\{x_{j} \mid 1 \leqslant j \leqslant n-1\right\}$ for some $n=n_{0}$. Obviously, there exists an infinite sequence of integers $\left\{n_{k}\right\}_{k=0}^{\infty}, n_{k+1}>n_{k}, k=0,1, \ldots$, such that $t$ is one of the $x_{j}$ whenever $n=n_{k}, k=0,1, \ldots$. With the exception of Theorems $1-3$, the notation $h \rightarrow 0$ will be assumed to mean that $n \rightarrow \infty$ through the sequence of integers $\left\{n_{k}\right\}_{k=0}^{\infty}$. We shall take $\sum_{j=m_{1}}^{\prime m_{2}} \alpha_{j}$ (or $\sum_{\substack{m_{2} \\ j=m_{1}}} \alpha_{j}$ ) to mean that $\alpha_{m_{2}}$ (or $\alpha_{m_{1}}$ ) is to be multiplied by $1 / 2$, while $\sum_{j=m_{1}}^{\prime m_{2}} \alpha_{j}$ will be taken to mean that both $\alpha_{m_{1}}$ and $\alpha_{m_{2}}$ are to be multiplied by $1 / 2$, in the respective summations.

Theorems 1-3, which are stated below without proof, form the basis of our development throughout the remainder of this work.

Theorem 1. Let the function $g(x)$ be $2 m$ times differentiable on $[a, b]$. Then

$$
\begin{align*}
D(h) & =\int_{a}^{b} g(x) d x-h \sum_{j=0}^{n} g\left(x_{j}\right) \\
& =\sum_{\mu=1}^{m-1} \frac{B_{2 \mu}}{(2 \mu)!}\left[g^{(2 \mu-1)}(a)-g^{(2 \mu-1)}(b)\right] h^{2 \mu}+R_{2 m}[g ;(a, b)] \tag{2.1}
\end{align*}
$$

where

$$
\begin{equation*}
R_{2 m}[g ;(a, b)]=h^{2 m} \int_{a}^{b} \frac{\bar{B}_{2 m}[(x-a) / h]-B_{2 m}}{(2 m)!} g^{(2 m)}(x) d x \tag{2.2}
\end{equation*}
$$

Here $B_{\mu}$ are the Bernoulli numbers, and $\bar{B}_{\mu}(x)$ is the periodic Bernoullian function of order $\mu$. In addition, since $\bar{B}_{\mu}(x)$ are bounded on $(-\infty, \infty)$, it follows that

$$
\begin{equation*}
\left|R_{2 m}[g ;(a, b)]\right| \leqslant M_{2 m}(b-a) h^{2 m} \max _{a \leqslant x \leqslant b}\left|g^{(2 m)}(x)\right| \tag{2.3}
\end{equation*}
$$

where

$$
\begin{equation*}
M_{2 m}=\max _{-\infty<x<\infty}\left|\bar{B}_{2 m}(x)-B_{2 m}\right| /(2 m)! \tag{2.4}
\end{equation*}
$$

and, therefore, is independent of $h$. Consequently, if $g(x)$ is infinitely differentiable on $[a, b]$, then $D(h)$ has an asymptotic expansion of the form

$$
\begin{equation*}
D(h) \sim \sum_{\mu=1}^{\infty} \frac{B_{2 \mu}}{(2 \mu)!}\left[g^{(2 \mu-1)}(a)-g^{(2 \mu-1)}(b)\right] h^{2 \mu} \quad \text { as } h \rightarrow 0 \tag{2.5}
\end{equation*}
$$

For a proof of this result see Steffensen (1950). The expansion in (2.5) is the classical Euler-Maclaurin expansion for trapezoidal rule approximations of integrals of smooth functions. Navot $(1961,1962)$ has extended the Euler Maxlaurin expansions to trapezoidal rule approximations of integrals of functions having algebraic and/or logarithmic end point singularities. By a different approach that utilizes generalized functions, Lyness and Ninham (1967) have rederived Navot's results. The results stated as Theorems 2 and 3 below are special cases of those proved by Navot.

Theorem 2. Let $g(x)$ be $2 m$ times differentiable on $[a, b]$ and let $G(x)=(x-a)^{s} g(x), s>-1$. Then

$$
\begin{align*}
D(h)= & \int_{a}^{b} G(x) d x-h \sum_{j=1}^{n} G\left(x_{j}\right) \\
= & -\sum_{\mu=1}^{m-1} \frac{B_{2 \mu}}{(2 \mu)!} G^{(2 \mu-1)}(b) h^{2 \mu} \\
& -\sum_{\mu=0}^{2 m-1} \frac{\zeta(-s-\mu)}{\mu!} g^{(\mu)}(a) h^{\mu+s+1}+\rho_{2 m} \tag{2.6}
\end{align*}
$$

where $\zeta(\tau)$ is the Riemann zeta function initially defined for $\operatorname{Re} \tau>1$ by $\zeta(\tau)=\sum_{k=1}^{\infty} k^{-\tau}$, and then continued analytically, and

$$
\begin{equation*}
\rho_{2 m}=O\left(h^{2 m}\right) \quad \text { as } h \rightarrow 0 \tag{2.7}
\end{equation*}
$$

If $g(x)$ is infinitely differentiable on $[a, b]$, then $D(h)$ has an asymptotic expansion of the form

$$
\begin{align*}
D(h) \sim & -\sum_{\mu=1}^{\infty} \frac{B_{2 \mu}}{(2 \mu)!} G^{(2 \mu-1)}(b) h^{2 \mu} \\
& -\sum_{\mu=0}^{\infty} \frac{\zeta(-s-\mu)}{\mu!} g^{(\mu)}(a) h^{\mu+s+1} \quad \text { as } h \rightarrow 0 \tag{2.8}
\end{align*}
$$

Starting from Theorem 2, Navot (1962) shows that the extensions of the Euler-Maclaurin expansion to trapezoidal rule approximations of integrals of function of the form $(x-a)^{s}[\log (x-a)]^{p} g(x)$, with $p$ being a
positive integer and $g(x)$ being sufficiently smooth in $[a, b]$, can be obtained by differentiating both sides of (2.6) $p$ times with respect to $s$. For $p=1$ the following result is obtained.

Theorem 3. Let $g(x)$ be $2 m$ times differentiable on [ $a, b$ ], and let $G(x)=(x-a)^{s} \log (x-a) g(x), s>-1$. Then

$$
\begin{align*}
D(h)= & \int_{a}^{b} G(x) d x-h \sum_{j=1}^{n} G\left(x_{j}\right) \\
= & -\sum_{\mu=1}^{m-1} \frac{B_{2 \mu}}{(2 \mu)!} G^{(2 \mu-1)}(b) h^{2 \mu} \\
& -\sum_{\mu=0}^{2 m-1}\left[-\zeta^{\prime}(-s-\mu)+\zeta(-s-\mu) \log h\right] \frac{g^{(\mu)}(a)}{\mu!} h^{\mu+s+1}+\tilde{\rho}_{2 m} \tag{2.9}
\end{align*}
$$

where $\zeta^{\prime}(\tau)=d \zeta(\tau) / d \tau$, and

$$
\begin{equation*}
\tilde{\rho}_{2 m}=0\left(h^{2 m}\right) \quad \text { as } h \rightarrow 0 \tag{2.10}
\end{equation*}
$$

If $g(x)$ is infinitely differentiable on $[a, b]$, then $D(h)$ has an asymptotic expansion of the form

$$
\begin{align*}
D(h) \sim & -\sum_{\mu=1}^{\infty} \frac{B_{2 \mu}}{(2 \mu)!} G^{(2 \mu-1)}(b) h^{2 \mu} \\
& -\sum_{\mu=0}^{\infty}\left[-\zeta^{\prime}(-s-\mu)+\zeta(-s-\mu) \log h\right] \frac{g^{(\mu)}(a)}{\mu!} h^{\mu+s+1} \tag{2.11}
\end{align*}
$$

Note that both Theorem 2 and Theorem 3 are true for any $s>-1$, although they were originally stated for $-1<s \leqslant 0$.

We shall now apply Theorems $1-3$ to integrals of the types

$$
\int_{a}^{b} \frac{g(x)}{x-t} d x \text { and } \int_{a}^{b}|x-t|^{s}(\log |x-t|)^{p} g(x) d x
$$

the former being defined as a Cauchy principal value integral.
Theorem 4. Let $g(x)$ be $2 m$ times differentiable on $[a, b]$, and let $G(x)=g(x) /(x-t)$. Then

$$
\begin{align*}
D(h)= & \int_{a}^{b} G(x) d x-h \sum_{\substack{j=0 \\
x_{j} \neq t}}^{n} G\left(x_{j}\right) \\
= & h g^{\prime}(t)+\sum_{\mu=1}^{m-1} \frac{B_{2 \mu}}{(2 \mu)!}\left[G^{(2 \mu-1)}(a)-G^{(2 \mu-1)}(b)\right] h^{2 \mu} \\
& +O\left(h^{2 m}\right) \quad \text { as } h \rightarrow 0 \tag{2.12}
\end{align*}
$$

Proof. We can assume without loss of generality that $t-a \leqslant b-t$. Then, since $t$ is one of the $x_{j}$, so is $b^{\prime}=2 t-a$. Furthermore, $t$ is the midpoint of the interval $\left[a, b^{\prime}\right]$. Now

$$
\begin{equation*}
\int_{a}^{b} G(x) d x=\int_{a}^{b^{\prime}} G(x) d x+\int_{b^{\prime}}^{b} G(x) d x \tag{2.13}
\end{equation*}
$$

and

$$
\begin{equation*}
\int_{a}^{b^{\prime}} G(x) d x=\int_{a}^{b^{\prime}} \frac{g(x)-g(t)}{x-t} d x \tag{2.14}
\end{equation*}
$$

the integral on the right-hand side of (2.14) being an ordinary integral, in which the integrand becomes $g^{\prime}(t)$ when $x=t$. Applying Theorem 1 to the right-hand side of (2.14), we have

$$
\begin{align*}
D_{1}(h)= & \int_{a}^{b^{\prime}} G(x) d x-h \sum_{\substack{x_{j} \leqslant b^{\prime} \\
x_{j} \neq t}}^{\prime \prime} \frac{g\left(x_{j}\right)-g(t)}{x_{j}-t}-h g^{\prime}(t) \\
= & \sum_{\mu=1}^{m-1} \frac{B_{2 \mu}}{(2 \mu)!}\left\{\left.\frac{d^{2 \mu-1}}{d x^{2 \mu-1}}\left[\frac{g(x)-g(t)}{x-t}\right]\right|_{x=a}\right. \\
& \left.-\left.\frac{d^{2 \mu-1}}{d x^{2 \mu-1}}\left[\frac{g(x)-g(t)}{x-t}\right]\right|_{x=b^{\prime}}\right\} h^{2 \mu} \\
& +O\left(h^{2 m}\right) \quad \text { as } h \rightarrow 0 \tag{2.15}
\end{align*}
$$

Now

$$
\begin{equation*}
\sum_{\substack{x_{j} \leqslant b^{\prime} \\ x_{j} \neq t}}^{\prime \prime} \frac{1}{x_{j}-t}=0 \tag{2.16}
\end{equation*}
$$

and

$$
\begin{equation*}
\frac{d^{r}}{d x^{r}}\left(\frac{1}{x-t}\right)=\frac{(-1)^{r} r!}{(x-t)^{r+1}}, \quad r=0,1, \ldots \tag{2.17}
\end{equation*}
$$

Combining (2.16) and (2.17) in (2.15), we have

$$
\begin{align*}
D_{1}(h) & =\int_{a}^{b^{\prime}} G(x) d x-h \sum_{\substack{x_{j} \leqslant b^{\prime} \\
x_{j} \neq t}}^{\prime \prime} G\left(x_{j}\right)-h g^{\prime}(t) \\
& =\sum_{\mu=1}^{m-1} \frac{B_{2 \mu}}{(2 \mu)!}\left[G^{(2 \mu-1)}(a)-G^{(2 \mu-1)}\left(b^{\prime}\right)\right] h^{2 \mu}+O\left(h^{2 m}\right) \quad \text { as } h \rightarrow 0 \tag{2.18}
\end{align*}
$$

Applying Theorem 1 to $\int_{b^{\prime}}^{b} G(x) d x$, we obtain

$$
\begin{align*}
D_{2}(h) & =\int_{b^{\prime}}^{b} G(x) d x-h \sum_{x_{j} \geqslant b^{\prime}}^{\prime \prime} G\left(x_{j}\right) \\
& =\sum_{\mu=1}^{m-1} \frac{B_{2 \mu}}{(2 \mu)!}\left[G^{(2 \mu-1)}\left(b^{\prime}\right)-G^{(2 \mu-1)}(b)\right] h^{2 \mu}+O\left(h^{2 m}\right) \quad \text { as } h \rightarrow 0 \tag{2.19}
\end{align*}
$$

Adding (2.18) and (2.19), (2.12) now follows.

Corollary. The remainder term $O\left(h^{2 m}\right)$ in (2.12) is actually given by

$$
\begin{equation*}
\tilde{R}_{2 m}[G ;(a, b)]=h^{2 m} \int_{a}^{b} \frac{\bar{B}_{2 m}[(x-a) / h]-B_{2 m}}{(2 m)!} G^{(2 m)}(x) d x \tag{2.20}
\end{equation*}
$$

this integral being interpreted as a Cauchy principal value integral.
Proof. The remainder term $O\left(h^{2 m}\right)$ in (2.15) is

$$
\begin{equation*}
R_{2 m}^{1}=h^{2 m} \int_{a}^{b^{\prime}} \frac{\bar{B}_{2 m}[(x-a) / h]-B_{2 m}}{(2 m)!} \frac{d^{2 m}}{d x^{2 m}}\left[\frac{g(x)-g(t)}{x-t}\right] d x \tag{2.21}
\end{equation*}
$$

Now making the change of variable $x=t+\xi$, we have

$$
\begin{align*}
& \int_{a}^{b^{\prime}} \frac{\bar{B}_{2 m}[(x-a) / h]-B_{2 m}}{(2 m)!} \frac{d^{2 m}}{d x^{2 m}}\left(\frac{1}{x-t}\right) d x \\
& \quad=\int_{-(t-a)}^{t-a} \frac{\bar{B}_{2 m}[\xi / h+(t-a) / h]-B_{2 m}}{(2 m)!} \frac{d^{2 m}}{d \xi^{2 m}}\left(\frac{1}{\xi}\right) d \xi \tag{2.22}
\end{align*}
$$

Recall that

$$
\begin{equation*}
\bar{B}_{k}(N+x)=\bar{B}_{k}(x), \quad N \text { integer }, \quad k \geqslant 2 \tag{2.23}
\end{equation*}
$$

and

$$
\begin{equation*}
\bar{B}_{k}(-x)=(-1)^{k} \bar{B}_{k}(x), \quad k \geqslant 2 \tag{2.24}
\end{equation*}
$$

Since $(t-a) / h$ is an integer, it follows using (2.23) and (2.24) that the integrand of the integral on the right-hand side of (2.22) is odd. Consequently, when taken as a Cauchy principal value integral, this integral is zero. The result now follows by using this in (2.21), and adding to $R_{2 m}^{1}$ the remainder term of the Euler-Maclaurin expansion for the integral $\int_{b^{\prime}}^{b} G(x) d x$.

Theorem 5. Let $g(x)$ be $2 m$ times differentiable on $[a, b]$, and let $G(x)=|x-t|^{s} g(x)$. Then

$$
\begin{align*}
D(h)= & \int_{a}^{b} G(x) d x-h \sum_{\substack{j=0 \\
x_{j} \neq t}}^{n} G\left(x_{j}\right) \\
= & \sum_{\mu=1}^{m-1} \frac{B_{2 \mu}}{(2 \mu)!}\left[G^{(2 \mu-1)}(a)-G^{(2 \mu-1)}(b)\right] h^{2 \mu} \\
& -2 \sum_{\mu=0}^{m-1} \frac{\zeta(-s-2 \mu)}{(2 \mu)!} g^{(2 \mu)}(t) h^{2 \mu+s+1}+O\left(h^{2 m}\right) \quad \text { as } h \rightarrow 0 \tag{2.25}
\end{align*}
$$

Proof. Applying Theorem 2 to $\int_{t}^{b} G(x) d x$, we have

$$
\begin{align*}
D_{1}(h)= & \int_{t}^{b} G(x) d x-h \sum_{x_{j}>t}^{\prime} G\left(x_{j}\right) \\
= & -\sum_{\mu=0}^{m-1} \frac{B_{2 \mu}}{(2 \mu)!} G^{(2 \mu-1)}(b) h^{2 \mu} \\
& -\sum_{\mu=0}^{2 m-1} \frac{\zeta(-s-\mu)}{\mu!} g^{(\mu)}(t) h^{\mu+s+1}+O\left(h^{2 m}\right) \quad \text { as } h \rightarrow 0 \tag{2.26}
\end{align*}
$$

Next applying Theorem 2 to $\int_{0}^{t-a} \xi^{s} g(t-\xi) d \xi\left[=\int_{a}^{t} G(x) d x\right]$, we have

$$
\begin{align*}
D_{2}(h)= & \int_{a}^{t} G(x) d x-h \sum_{x_{j}<t}^{*} G\left(x_{j}\right) \\
= & \sum_{\mu=1}^{m-1} \frac{B_{2 \mu}}{(2 \mu)!} G^{(2 \mu-1)}(a) h^{2 \mu} \\
& -\sum_{\mu=0}^{2 m-1}(-1)^{\mu} \frac{\zeta(-s-\mu)}{\mu!} g^{(\mu)}(t) h^{\mu+s+1}+O\left(h^{2 m}\right) \quad \text { as } h \rightarrow 0 \tag{2.27}
\end{align*}
$$

Adding (2.26) and (2.27), (2.25) follows.

Theorem 6. Let $G(x)=|x-t|^{s} \log |x-t| g(x)$ in the statement of Theorem 5, everything else being the same. Then

$$
\begin{align*}
D(h)= & \int_{a}^{b} G(x) d x-h \sum_{\substack{j=0 \\
x_{j} \neq t}}^{n} G\left(x_{j}\right) \\
= & \sum_{\mu=1}^{m-1} \frac{B_{2 \mu}}{(2 \mu)!}\left[G^{(2 \mu-1)}(a)-G^{(2 \mu-1)}(b)\right] h^{2 \mu} \\
& -2 \sum_{\mu=0}^{m-1}\left[-\zeta^{\prime}(-s-2 \mu)+\zeta(-s-2 \mu) \log h\right] \\
& \times \frac{g^{(2 \mu)}(t)}{(2 \mu)!} h^{2 \mu+s+1}+O\left(h^{2 m}\right) \quad \text { as } h \rightarrow 0 \tag{2.28}
\end{align*}
$$

Proof. Similar to that of Theorem 5.
Corollary. For $s=0$, (2.28) becomes

$$
\begin{align*}
D(h)= & g(t) h \log h+\sum_{\mu=1}^{m-1} \frac{B_{2 \mu}}{(2 \mu)!}\left[G^{(2 \mu-1)}(a)-G^{(2 \mu-1)}(b)\right] h^{2 \mu} \\
& +2 \sum_{\mu=0}^{m-1} \frac{\zeta^{\prime}(-2 \mu)}{(2 \mu)!} g^{(2 \mu)}(t) h^{2 \mu+1}+O\left(h^{2 m}\right) \quad \text { as } h \rightarrow 0 \tag{2.29}
\end{align*}
$$

Proof. The proof follows by setting $s=0$ in (2.28) and using the facts that

$$
\begin{equation*}
\zeta(0)=-1 / 2, \quad \zeta(-2 \mu)=0, \quad \mu=1,2, \ldots \tag{2.30}
\end{equation*}
$$

(see Abramowitz and Stegun, 1964, p. 807).
The results in the following theorem will be the ones on which our quadrature methods will be based.

Theorem 7. Assume that the functions $g(x)$ and $\tilde{g}(x)$ are $2 m$ times differentiable on $[a, b]$. Assume also that the functions $G(x)$ are periodic with period $T=b-a$, and that they are $2 m$ times differentiable on $\tilde{R}=(-\infty, \infty) \backslash\{t+k T\}_{k=-\infty}^{\infty}$. Then we have the following:
(a) If $G(x)=g(x) /(x-t)+\tilde{g}(x)$, and

$$
\begin{equation*}
Q_{n}[G]=h \sum_{\substack{j=1 \\ x_{j} \neq i}}^{n} G\left(x_{j}\right) \tag{2.31a}
\end{equation*}
$$

then

$$
\begin{equation*}
E_{n}[G]=\left[\tilde{g}(t)+g^{\prime}(t)\right] h+O\left(h^{2 m}\right) \quad \text { as } h \rightarrow 0 \tag{2.32a}
\end{equation*}
$$

(b) If $G(x)=|x-t|^{s} g(x)+\tilde{g}(x), s>-1$, and

$$
\begin{equation*}
Q_{n}[G]=h \sum_{\substack{j=1 \\ x_{j} \neq i}}^{n} G\left(x_{j}\right)+\tilde{g}(t) h-2 \zeta(-s) g(t) h^{s+1} \tag{2.31b}
\end{equation*}
$$

then

$$
\begin{equation*}
E_{n}[G]=-2 \sum_{\mu=1}^{m-1} \frac{\zeta(-s-2 \mu)}{(2 \mu)!} g^{(2 \mu)}(t) h^{2 \mu+s+1}+O\left(h^{2 m}\right) \quad \text { as } h \rightarrow 0 \tag{2.32b}
\end{equation*}
$$

(c) If $G(x)=|x-t|^{s} \log |x-t| g(x)+\tilde{g}(x), s>-1$, and

$$
\begin{equation*}
Q_{n}[G]=h \sum_{\substack{j=1 \\ x_{j} \neq t}}^{n} G\left(x_{j}\right)+\tilde{g}(t) h+2\left[\zeta^{\prime}(-s)-\zeta(-s) \log h\right] g(t) h^{s+1} \tag{2.31c}
\end{equation*}
$$

then

$$
\begin{align*}
E_{n}[G]= & 2 \sum_{\mu=1}^{m-1}\left[\zeta^{\prime}(-s-2 \mu)-\zeta(-s-2 \mu) \log h\right] \frac{g^{(2 \mu)}(t)}{(2 \mu)!} \\
& \times h^{2 \mu+s+1}+O\left(h^{2 m}\right) \quad \text { as } h \rightarrow 0 \tag{2.32c}
\end{align*}
$$

( $\mathrm{c}^{\prime}$ ) When $s=0$ in (c), by (2.30) and $\zeta^{\prime}(0)=-\frac{1}{2} \log (2 \pi)$, (see Abramowitz and Stegun, 1964, p. 807), (2.31c) and (2.32c) reduce to

$$
Q_{n}[G]=h \sum_{\substack{j=1 \\ x_{j} \neq t}}^{n} G\left(x_{j}\right)+\tilde{g}(t) h+\log \left(\frac{h}{2 \pi}\right) g(t) h
$$

and

$$
E_{n}[G]=2 \sum_{\mu=1}^{m-1} \frac{\zeta^{\prime}(-2 \mu)}{(2 \mu)!} g^{(2 \mu)}(t) h^{2 \mu+1}+O\left(h^{2 m}\right) \quad \text { as } h \rightarrow 0
$$

where $E_{n}[G]=\int_{a}^{b} G(x) d x-Q_{n}[G]$ in all cases.
Remarks. (1) As is seen from (2.32b), (2.32c), and (2.32c'), $E_{n}[G]$ depends only on $t$, and is independent of $a$ and $b$. This is a consequence of the periodicity of $G(x)$, and is an important property that will be exploited in the derivation of Romberg-type quadrature formulas in the next section.
(2) Until now we assumed that $t$ is one of the points $x_{j}$. When $t$ is arbitrary the periodicity of $G(x)$ can be used to shift the interval $[a, b]$ to $\left[a^{\prime}, b^{\prime}\right]$ such that $t$ coincides with one of the $x_{j}$ in the new interval $\left[a^{\prime}, b^{\prime}\right]$.

This, combined with the observation of Remark 1, means that (a) $E_{n}[G]$ in all parts of Theorem 7 stays the same if the sum $\sum_{j=1, x_{j} \neq t}^{n} G\left(x_{j}\right)$ in (2.31a), (2.31b), (2.31c), and (2.31c $)$ is replaced by $\sum_{j=1}^{n-1} G(t+j h)$ or by the identical sum $\sum_{j \neq 0, a<t+j h \leqslant b} G(t+j h)$, and (b) Theorem 7 holds for all positive integers $n$. These facts will be repeatedly used in Section 4 without further explanation.
(3) In each of the cases of Theorem 7, the numerical quadrature formula $Q_{n}[G]$ is computed by using the function values only. The quadrature formula (2.31a) has an error that is of order $h$, and it would seem that one would have to know $g^{\prime}(t)$ with a high accuracy in order to improve this formula. But as we shall see in the next section, the term $\left[\tilde{g}(t)+g^{\prime}(t)\right] h$ that appears in $E_{n}[G]$ is easily removed by one extrapolation.
(4) Let $G(x)=\sum_{k=1}^{M} g_{k}(x)|x-t|^{\alpha_{k}}(\log |x-t|)^{p_{k}}+\hat{g}_{1}(x) /(x-t)+\hat{g}_{2}(x)$, where $\alpha_{k}$ and $p_{k}$ are as described following (1.3), and $g_{k}(x), 1 \leqslant k \leqslant M$, $\hat{g}_{1}(x)$, and $\hat{g}_{2}(x)$ are $2 m$ times differentiable on $[a, b]$ and $G(x)$ is periodic with period $T=b-a$ and is $2 m$ times differentiable on $\widetilde{R}$. Inspection of Theorems 4-7 reveals that if we form $Q_{n}[G]$ as the sum of the quadrature formulas for each one of the terms in $G(x)$, then the error $E_{n}[G]$ does not contain any contribution from $G(x)$ and its derivatives at the end points, and the only contribution to $E_{n}[G]$ comes from $G(x)$ and its derivatives at $x=t$, as in Theorem 7 .

## 3. ROMBERG-TYPE NUMERICAL QUADRATURE FORMULAS

Using the results of Theorem 7, we can apply extrapolation techniques to derive Romberg-type numerical quadrature formulas for $\int_{a}^{b} G(x) d x$. The simplest case is that of Theorem 7a, and we deal with it first.

Theorem 8. Let $G(x)$ and $Q_{n}[G]$ be as in Theorem 7a. Let $h_{k}=T / k$ and

$$
\begin{equation*}
\widetilde{Q}_{n}[G]=2 Q_{2 n}[G]-Q_{n}[G]=h_{n} \sum_{j=1}^{n} G\left(a+j h_{n}-h_{n} / 2\right) \tag{3.1}
\end{equation*}
$$

i.e., $\widetilde{Q}_{n}[G]$ is a midpoint rule approximation. Then

$$
\begin{equation*}
\tilde{E}_{n}[G]=\int_{a}^{b} G(x) d x-\tilde{Q}_{n}[G]=O\left(h_{n}^{2 m}\right) \quad \text { as } h_{n} \rightarrow 0 \tag{3.2}
\end{equation*}
$$

Proof. Equation (3.2) follows directly from Theorem 7a, the $h$ term in $E_{n}[G]$ being eliminated when $\widetilde{Q}_{n}[G]$ is formed.

As a result of Theorem 8 we conclude that if $G(x)$ is infinitely differentiable on $\widetilde{R}=(-\infty, \infty) \backslash\{t+k T\}_{k=-\infty}^{\infty}$, then $\tilde{E}_{n}[G]$ tends to zero more quickly than any inverse power of $n$, as $n \rightarrow \infty$. We can improve this result considerably whenever $G(z)$ is analytic in a strip in the complex $z$ plane, which, with the exception of the points $t+k T, k=0, \pm 1, \ldots$, contains the real line $\operatorname{Im} z=y=0$ in its interior.

Theorem 9. Let $G(z)$ in the previous theorem be analytic in the strip $|\operatorname{Im} z|<\sigma$, except at the simple poles $t+k T, k=0, \pm 1, \pm 2, \ldots$. Then

$$
\begin{equation*}
\left|\tilde{E}_{n}[G]\right| \leqslant 2 T M\left(\sigma^{\prime}\right) \frac{\exp \left[-2 \pi n \sigma^{\prime} / T\right]}{1-\exp \left[-2 \pi n \sigma^{\prime} / T\right]}, \quad \sigma^{\prime}<\sigma \tag{3.3}
\end{equation*}
$$

where

$$
\begin{equation*}
M(\tau)=\max \left\{\max _{-\infty<x<\infty}\left|G_{e}(x+i \tau)\right|, \max _{-\infty<x<\infty}\left|G_{e}(x-i \tau)\right|\right\} \tag{3.4}
\end{equation*}
$$

and

$$
\begin{equation*}
G_{e}(\xi)=\frac{1}{2}[G(t+\zeta)+G(t-\xi)] \tag{3.5}
\end{equation*}
$$

Proof. By the periodicity of $G(x)$, we can write

$$
\begin{align*}
\int_{a}^{b} G(x) d x & =\int_{t-T / 2}^{t+T / 2} G(x) d x \\
& =\int_{-T / 2}^{T / 2} G_{e}(\xi) d \xi \tag{3.6}
\end{align*}
$$

$G_{e}(z)$ is analytic in the strip $|\operatorname{Im} z|<\sigma$ and is periodic with period $T$. After some algebra it can be shown that the $n$ (odd)-point trapezoidal rule approximation or the $n$ (even)-point midpoint rule approximation to $\int_{-T / 2}^{T / 2} G_{e}(\xi) d \xi$ is just $\widetilde{Q}_{n}[G]$. Applying now a theorem due to Davis (1955) (see also Davis and Rabinowitz, 1984, pp.314-316), (3.3)-(3.5) follow. (Actually Davis' result is stated for the trapezoidal rule. However, inspection of his proof shows it to be valid for the off-set trapezoidal rule. The midpoint rule is an off-set trapezoidal rule.)

For arbitrary $t$, by Remark 2 following Theorem 7, the approximation $\tilde{Q}_{n}[G]$ can be replaced by

$$
\begin{equation*}
\widetilde{Q}_{n}[G]=h_{n} \sum_{j=1}^{n} G\left(t+j h_{n}-h_{n} / 2\right) \tag{3.7}
\end{equation*}
$$

and again Theorems 8 and 9 apply.
$\widetilde{Q}_{n}[G]$ in Theorem 8 has been obtained by employing the Richardson extrapolation process once, eliminating the term $O(h)$ in the error $E_{n}[G]$.

Since $\tilde{E}_{n}[G]=O\left(h^{2 m}\right)$ with $m$ as large as we wish, there is no need for further extrapolation. For $Q_{n}[G]$ as given in Theorem $7 \mathrm{~b}, \mathrm{c}, \mathrm{c}^{\prime}$, however, we apply the Richardson extrapolation process (or generalizations of it) repeatedly in order to eliminate successive terms in $E_{n}[G]$, thus obtaining Romberg-type numerical quadrature formulas with higher degrees of accuracy. (A summary of those extrapolation methods relevant to the present work is given in the appendix.)

Let $G(x)$ and $Q_{n}[G]$ be as in Theorem $7 \mathrm{~b}, \mathrm{c}, \mathrm{c}^{\prime}$, and define $A=\int_{a}^{b} G(x) d x$ and $A(h)=Q_{n}[G]$ in the notation of the appendix. Select a sequence of integers $\left\{n_{l}\right\}_{l=0}^{\infty}, 1 \leqslant n_{0}<n_{1}<\cdots$, and set $h_{l}=T / n_{i}, l=0,1, \ldots$. Obviously $\lim _{l \rightarrow \infty} h_{l}=0$, as required; see the appendix. Now the Rombergtype quadrature formulas $A_{q}^{(m)}$, based on $A\left(h_{l}\right), m \leqslant l \leqslant m+q$, are of the form

$$
\begin{equation*}
A_{q}^{(m)}=\sum_{k=0}^{q} d_{q, k}^{(m)} A\left(h_{m+k}\right) \tag{3.8}
\end{equation*}
$$

where $d_{q, k}^{(m)}, 0 \leqslant k \leqslant q$, are constants determined by the nature of the asymptotic expansion of $A-A(h)=E_{n}[G]$ as $h \rightarrow 0$, as described in the appendix. Some of the details for two special cases are given below:
(1) If $G(x)$ is as in Theorem 7 b (or Theorem $7 \mathrm{c}^{\prime}$ ), then $A(h)$ is of the form given in (a) of the appendix with $\gamma_{i}=s+1+2 i$, and $\beta_{i}=$ $-2 \zeta(-s-2 i) g^{(2 i)}(t) /(2 i)!$ [or $\gamma_{i}=1+2 i$, and $\left.\beta_{i}=2 \zeta^{\prime}(-2 i) g^{(2 i)}(t) /(2 i)!\right]$, $i=1,2, \ldots$. Hence for a sequence of the form $h_{l}=h_{0} \rho^{l}, l=0,1, \ldots$, the $d_{n, k}^{(j)}$ can be computed from (A.6). For arbitrary $h_{l}$, the algorithm given in (A.11) and (A.12) in (b) of the appendix is appropriate with $\varphi(h)=h^{s+3}$ and $r=2$ [or $\varphi(h)=h^{3}$ and $r=2$ ].
(2) If $G(x)$ is as in Theorem 7c with $s \neq 0$, then $A(h)$ is of the form given in (A.1) with $e_{i}(h)=\left[\zeta^{\prime}(-s-2 i)-\zeta(-s-2 i) \log h\right] h^{2 i+s+1}$, and $\beta_{i}=2 g^{(2 i)}(t) /(2 i)!, i=1,2, \ldots$ The $d_{n, k}^{(j)}$ then can be obtained by solving the equations in (A.5).

Before closing this section we recall that, for any positive integer $n, A(h)=Q_{n}[G]$ in Theorem 7b, c, $\mathrm{c}^{\prime}$ has the form

$$
\begin{equation*}
A(h)=h \sum_{\substack{j \neq 0 \\ a<t+j h \leqslant b}} G\left(x_{j}\right)+C(t, h) \tag{3.9}
\end{equation*}
$$

where

$$
\tilde{g}(t) h-2 \zeta(-s) g(t) h^{s+1} \quad \text { for Theorem } 7 \mathrm{~b}
$$

$$
\begin{align*}
C(t, h)= & \tilde{g}(t) h+2\left[\zeta^{\prime}(-s)-\zeta(-s) \log h\right] g(t) h^{s+1} \quad \text { for Theorem } 7 \mathrm{c} \\
& \tilde{g}(t) h+\log \left(\frac{h}{2 \pi}\right) g(t) h \quad \text { for Theorem } 7 \mathrm{c}^{\prime} \tag{3.10}
\end{align*}
$$

## 4. THE QUADRATURE METHODS FOR INTEGRAL EQUATIONS

In what follows we consider the integral equation (1.1), with $f(t)$ and $g(t)$ being periodic in $t$ with period $T$, and $K(t, x)$ being periodic both in $t$ and $x$ with period $T$. Of course, a similar treatment can be given to the eigenvalue problem in (1.2).

### 4.1. The Singular Case

Let $K(t, x)=\hat{H}_{1}(t, x) /(t-x)+\hat{H}_{2}(t, x)$. For a given integer $N$, let $h \equiv h_{2 N}=(b-a) /(2 N)$, and $x_{j}=a+j h, j=1, \ldots, 2 N$. Then setting $t=x_{i}$ for some $i$, and approximating the integral $\int_{a}^{b} K\left(x_{i}, x\right) f(x) d x$ by the rule $\widetilde{Q}_{N}$ in (3.1), we write down the following set of equations for the $2 N$ unknowns $\tilde{f}_{j}$ [the approximations to the corresponding $\left.f\left(x_{j}\right)\right]$ :

$$
\begin{equation*}
\omega \widetilde{f}_{i}+2 h \sum_{j=1}^{2 N} \varepsilon_{i j} K\left(x_{i}, x_{j}\right) \widetilde{f}_{j}=g\left(x_{i}\right), \quad i=1,2, \ldots, 2 N \tag{4.1}
\end{equation*}
$$

where

$$
\varepsilon_{i j}=\left\{\begin{array}{lll}
1 & \text { if } & |i-j| \text { odd }  \tag{4.2}\\
0 & \text { if } & |i-j| \text { even }
\end{array}\right.
$$

### 4.2. The Weakly Singular Case

We mentioned in the previous section that when $G(x)$ is a known function any set of integers $\left\{n_{l}\right\}_{l=0}^{\infty}, 1 \leqslant n_{0}<n_{1}<\cdots$, can be chosen for computing the approximation $A_{q}^{(m)}$ to the integrals $\int_{a}^{b} G(x) d x$ in Theorem $7 \mathrm{~b}, \mathrm{c}, \mathrm{c}^{\prime}$. If, however, we want to use the Romberg-type formula $A_{q}^{(m)}$ for solving integral equations by quadrature methods, the $n_{l}$ cannot be arbitrary. In fact, we should choose the $n_{l}$ (hence the $h_{l}=T / n_{l}$ ), such that the sets of abscissas that enter the computation of $A\left(h_{m+k}\right)=Q_{n_{m+k}}[G]$, $0 \leqslant k \leqslant q-1$, where $G(x)=K(t, x) f(x)$, are all subsets of the set of abscissas that enter the computation of $A\left(h_{m+q}\right)$. This is achieved by picking the $n_{l}, m \leqslant l \leqslant m+q-1$, as divisors of $n_{m+q}$. With this choice of the $n_{l}$ let $x_{j}=a+j h_{m+q}=a+j T / n_{m+q}, 1 \leqslant j \leqslant n_{m+q}$. With the help of (3.9) it can be verified that, for $t=x_{i}$,

$$
\begin{equation*}
Q_{n_{m+k}}[G]=h_{m+k} \sum_{\substack{j=1 \\ j \neq i}}^{n_{m+q}} \varepsilon_{i j}^{m, q, k} G\left(x_{j}\right)+C\left(x_{i}, h_{m+k}\right) \tag{4.3}
\end{equation*}
$$

where

$$
\varepsilon_{i j}^{m, q, k}= \begin{cases}1 & \text { if }|i-j| \text { is divisible by } n_{m+q} / n_{m+k}  \tag{4.4}\\ 0 & \text { otherwise }\end{cases}
$$

Thus (3.8) becomes, for $t=x_{i}$,

$$
\begin{equation*}
A_{q}^{(m)}=\sum_{\substack{j=1 \\ j \neq i}}^{n_{m+q}}\left(\sum_{k=0}^{q} \varepsilon_{i j}^{m, q, k} d_{q, k}^{(m)} h_{m+k}\right) G\left(x_{j}\right)+\sum_{k=0}^{q} d_{q, k}^{(m)} C\left(x_{i}, h_{m+k}\right) \tag{4.5}
\end{equation*}
$$

Now for $j \neq i, G\left(x_{j}\right)=K\left(x_{i}, x_{j}\right) f\left(x_{j}\right)$. By Theorem $7 \mathrm{~b}, \mathrm{c}, \mathrm{c}^{\prime}$ we note that when
(b) $K(t, x)=H_{1}(t, x)|t-x|^{s}+\hat{H}_{2}(t, x), g(x)=H_{1}(t, x) f(x)$ and $\tilde{g}(x)=$ $\hat{H}_{2}(t, x) f(x)$ in (2.31b). Thus $C(t, h)=\hat{C}(t, h) f(t)$, where

$$
\begin{equation*}
\hat{C}(t, h)=h\left[\hat{H}_{2}(t, t)-2 \zeta(-s) H_{1}(t, t) h^{s}\right] \tag{4.6b}
\end{equation*}
$$

(c) $K(t, x)=H_{1}(t, x)|t-x|^{s} \log |t-x|+\hat{H}_{2}(t, x), g(x)$ and $\tilde{g}(x)$ in (2.31c) are as in (b) above. Thus $C(t, h)=\hat{C}(t, h) f(t)$, where

$$
\begin{equation*}
\hat{C}(t, h)=h\left\{\hat{H}_{2}(t, t)+2\left[\zeta^{\prime}(-s)-\zeta(-s) \log h\right] H_{1}(t, t) h^{s}\right\} \tag{4.6c}
\end{equation*}
$$

(c') $K(t, x)=H_{1}(t, x) \log |t-x|+\hat{H}_{2}(t, x), g(x)$ and $\tilde{g}(x)$ in (2.31c $)$ are as in (b) and (c) above. Thus $C(t, h)=\hat{C}(t, h) f(t)$, where

$$
\hat{C}(t, h)=h\left[\hat{H}_{2}(t, t)+\log \left(\frac{h}{2 \pi}\right) H_{1}(t, t)\right]
$$

Combining the above in (4.5), approximating the integral $\int_{a}^{b} K(t, x) f(x) d x$ in (1.1) by $A_{q}^{(m)}$, and replacing the $f\left(x_{j}\right)$ by the corresponding approximations $f_{j}, 1 \leqslant j \leqslant n_{m+q}$, we obtain the appropriate quadrature methods for (1.1), which are defined by the systems of linear equations

$$
\begin{equation*}
\omega \tilde{f}_{i}+\sum_{j=1}^{n_{m+q}} \tilde{K}_{i j} \tilde{f}_{j}=g\left(x_{i}\right), \quad 1 \leqslant i \leqslant n_{m+q} \tag{4.7}
\end{equation*}
$$

where

$$
\begin{equation*}
\tilde{K}_{i j}=\left(\sum_{k=0}^{q} \varepsilon_{i j}^{m, q, k} d_{q, k}^{(m)} h_{m+k}\right) K\left(x_{i}, x_{j}\right), \quad j \neq i \tag{4.8}
\end{equation*}
$$

and

$$
\begin{equation*}
\tilde{K}_{i i}=\sum_{k=0}^{q} d_{q, k}^{(m)} \hat{C}\left(x_{i}, h_{m+k}\right) \tag{4.9}
\end{equation*}
$$

with $\hat{C}(t, h)$ as defined in (4.6b), (4.6c), and (4.6c').
It is not the purpose of this work to give precise error bounds or convergence results for $\Lambda=\max _{j}\left|f\left(x_{j}\right)-\tilde{f}_{j}\right|$. However, in general, we would
expect $A$ to be of the order of magnitude of the error in the numerical quadrature formula used in approximating the integral $\int_{a}^{b} K(t, x) f(x) d x$. Thus, for the singular case, if $K(t, z)$ is meromorphic in the strip $|\operatorname{Im} z|<\sigma$ with its only poles at $t+k T, k=0, \pm 1, \pm 2, \ldots$, and $g(z)$ is analytic in the same strip, we would expect $f(z)$ to be analytic in this strip too. Therefore, by Theorem $9, \int_{a}^{b} G(x) d x-\tilde{Q}_{n}[G]=O\left(e^{-2 \pi N \sigma / T}\right)$ as $N \rightarrow \infty$, thus we would expect $A=O\left(e^{-2 \pi N \sigma / T}\right)$ as $N \rightarrow \infty$ too. For the weakly singular case, by the appendix, $\int_{a}^{b} G(x) d x-A_{q}^{(m)}=O\left(e_{q+1}\left(h_{m}\right)\right)$ as $m \rightarrow \infty$, in general. Thus, for $M=1, s=\alpha_{1}$, and $p=p_{1}$ in (1.3), we would expect $A=O\left(h_{m}^{s+3+2 q}(\log h)^{p}\right)$ as $m \rightarrow \infty$ for $s \neq 0, p=0,1$, and $A=O\left(h_{m}^{3+2 q}\right)$ as $m \rightarrow \infty$ for $s=0, p=1$.

Finally, we could use the approximations $\tilde{f}_{i}$ to $f\left(x_{i}\right), i=1, \ldots, \hat{N}$, where $\hat{N}=2 N$, for singular equations and $\hat{N}=n_{m+q}$ for weakly singular equations, to construct a trigonometric interpolation polynomial $P_{\hat{N}}(x)$ in $\cos (2 \pi k x / T), \sin (2 \pi k x / T), k=0,1, \ldots$, satisfying $P_{\hat{N}}\left(x_{i}\right)=\widetilde{f}_{i}, i=1, \ldots, \hat{N}$, thus obtaining an approximation to $f(x)$ for all $x$ in $[a, b]$.

## 5. NUMERICAL EXAMPLES

### 5.1. The Singular Case

Example 1 (See Mikhlin, 1964, pp. 122-124).

$$
\begin{equation*}
a f(t)+\frac{b}{2 \pi} \int_{0}^{2 \pi} \cot \left(\frac{x-t}{2}\right) f(x) d x=u(t) \tag{5.1}
\end{equation*}
$$

When $a \neq 0$ and $a^{2}+b^{2} \neq 0(a$ and $b$ may be complex) a unique solution exists and is given by

$$
\begin{align*}
f(t)= & \frac{a}{a^{2}+b^{2}} u(t)-\frac{b}{2 \pi\left(a^{2}+b^{2}\right)} \int_{0}^{2 \pi} u(x) \cot \left(\frac{x-t}{2}\right) d x \\
& +\frac{b^{2}}{2 \pi a\left(a^{2}+b^{2}\right)} \int_{0}^{2 \pi} u(x) d x \tag{5.2}
\end{align*}
$$

We first observe that the kernel function $K(t, x)=(b / 2 \pi) \cot [(x-t) / 2]$ is periodic with period $2 \pi$ in both $x$ and $t$. Also for fixed $t, K(t, z)$ is meromorphic in the whole $z$ plane with simple poles at $t+2 \pi k$, $k=0, \pm 1, \pm 2, \ldots$, thus being of the form described in Section 4.1. Next, we observe that if $u(t)$ is periodic with period $2 \pi$, then so is $f(t)$. Also if $u(z)$ is analytic in a strip $|\operatorname{Im} z|<\sigma$, then so is $f(z)$. The last two assertions can be verified with the help of (5.2).

In our numerical experiments we chose $u(t)=(D+\cos t)^{-1}, D>1$, so that both $u(z)$ and $f(z)$ are analytic in the strip $|\operatorname{Im} z|<\sigma=$ $\log \left[D+\left(D^{2}-1\right)^{1 / 2}\right]$. With this choice of $u(t),(5.2)$ becomes

$$
f(t)=\frac{1}{a^{2}+b^{2}}\left\{\frac{1}{D+\cos t}\left(a-\frac{b \sin t}{\left(D^{2}-1\right)^{1 / 2}}\right)+\frac{b^{2}}{a\left(D^{2}-1\right)^{1 / 2}}\right\}
$$

Denote $2 N$, the number of abscissas in (4.1), by $\hat{N}$. Then $h=2 \pi / \hat{N}$ and $x_{j}=j h, \quad 1 \leqslant j \leqslant \hat{N}$. Denote $\vec{f}_{\hat{N}, j} \equiv \vec{f}_{j}, \quad 1 \leqslant j \leqslant \hat{N}, \quad$ and let $\Lambda_{\hat{N}}=$ $\max _{1 \leqslant j \leqslant \hat{N}}\left|f\left(x_{j}\right)-\tilde{f}_{\hat{N}, j}\right|$. Then, by what has been said in the paragraph following (4.9), we would expect to have $\Lambda_{\hat{N}}=O\left(e^{-\sigma \hat{N} / 2}\right)$ as $\hat{N} \rightarrow \infty$. This is born out by the numerical results, where, for each $\hat{N}, \sigma$ can be estimated by the formula

$$
\begin{equation*}
\sigma \approx \frac{2}{\hat{N}_{\max }-\hat{N}} \log \frac{A_{\hat{N}}}{A_{\hat{N}_{\max }}} \equiv \sigma_{\hat{N}, \hat{N}_{\max }} \tag{5.3}
\end{equation*}
$$

where $\hat{N}_{\text {max }}$ is the maximum of the $\hat{N}^{\prime} s$ used in the computation. This estimate, of course, is based on the $\Lambda_{\hat{N}}$, which in turn are based on $f(x)$. An estimate based solely on the computed values $\tilde{f}_{\hat{N}, j}$ can be obtained from the formula

$$
\begin{equation*}
\sigma \approx \frac{2}{J} \log \left|\frac{\tilde{f}_{\hat{N}+2 J, \hat{N}+2 J}-\widetilde{f}_{\hat{N}+J, \hat{N}+J}}{\widetilde{f}_{\hat{N}+J, \hat{N}+J}-\widetilde{f}_{\hat{N}, \hat{N}}}\right| \equiv \sigma_{\hat{N}} \tag{5.4}
\end{equation*}
$$

This follows from the following expected behavior of $\hat{f}_{\hat{N}, \hat{N}}$ :

$$
\begin{equation*}
\tilde{f}_{\hat{N}, \hat{N}} \sim f\left(x_{\hat{N}}\right)+C\left(x_{\hat{N}}\right) e^{-\sigma \hat{N} / 2} \quad \text { as } \quad \hat{N} \rightarrow \infty \tag{5.5}
\end{equation*}
$$

Here, of course, $\tilde{f}_{\hat{N}, \hat{N}}$ and $\hat{x}_{N}=2 \pi$ can be replaced by $\tilde{f}_{\hat{N}, j(\hat{N})}$ and $\bar{x}$, respectively, where $x_{j(\hat{N})}=\bar{x}$ is the same for all $\hat{N}$ used in the computations [in (5.4) they are $\hat{N}, \hat{N}+J$, and $\hat{N}+2 J]$.

Table I gives the results obtained for $\Lambda_{\hat{N}}, \sigma_{\hat{N}, \hat{N}_{\max }}$, and $\sigma_{\hat{N}}, \hat{N}=4(4) 44$, $\hat{N}_{\max }=44$, with $a=b=1$. Note that $\sigma_{44,44}, \sigma_{4}$, and $\sigma_{12}$ are not defined. Note also that, for $D=2, \sigma_{\hat{N}, \hat{N}_{\max }}$ and $\sigma_{\hat{N}}$ deteriorate for $\hat{N}$ large. This is due to the fact that there is a loss of significance in the arguments of $\log$ in (5.3) and (5.4), which is caused by the high accuracy of the $\tilde{f}_{\hat{N}, j}$.

### 5.2. The Weakly Singular Case

In both of the examples below the kernel function is of the form described in ( $\mathrm{c}^{\prime}$ ) of Section 4.2, namely, $K(t, x)=H_{1}(t, x) \times$ $\log |t-x|+\hat{H}_{2}(t, x)$. Therefore, the approximations $\hat{f}_{j}$ to $f\left(x_{j}\right)$, the solution

Table I. Results for $\boldsymbol{A}_{\hat{N}}, \sigma_{\hat{N}, \hat{N}_{\text {max }}}$, and $\sigma_{\hat{N}}, \hat{N}=4(4) 44, \hat{N}_{\text {max }}=44$, for Example $1^{a}$

| $\hat{N}$ | $D=1.1$ |  |  | $D=2$ |  |  |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: |
|  | $A_{N}$ | $\sigma_{\hat{N}, \hat{N}_{\text {max }}}$ | $\sigma_{\hat{N}}$ | $\Lambda_{\hat{N}}$ | $\sigma_{\hat{N}, \hat{N}_{\text {max }}}$ | $\sigma_{\hat{N}}$ |
| 4 | $2.03 \times 10^{0}$ | 0.426 |  | $6.10 \times 10^{-2}$ | 1.3136 |  |
| 8 | $1.12 \times 10^{0}$ | 0.4407 |  | $4.60 \times 10^{-3}$ | 1.3157 |  |
| 12 | $4.93 \times 10^{-1}$ | 0.445 | 0.67 | $3.37 \times 10^{-4}$ | 1.3170 | 1.354 |
| 16 | $2.01 \times 10^{-1}$ | 0.4442 | 0.52 | $2.41 \times 10^{-5}$ | 1.31677 | 1.3195 |
| 20 | $7.98 \times 10^{-2}$ | 0.4411 | 0.474 | $1.73 \times 10^{-6}$ | 1.31683 | 1.3171 |
| 24 | $3.33 \times 10^{-2}$ | 0.4420 | 0.456 | $1.25 \times 10^{-7}$ | 1.3171 | 1.316971 |
| 28 | $1.39 \times 10^{-2}$ | 0.44339 | 0.4485 | $8.94 \times 10^{-9}$ | 1.3169586 | 1.3169588 |
| 32 | $5.73 \times 10^{-3}$ | 0.44303 | 0.4456 | $6.42 \times 10^{-10}$ | 1.316996 | 1.3169580 |
| 36 | $2.33 \times 10^{-3}$ | 0.4400 | 0.4444 | $4.62 \times 10^{-11}$ | 1.3173 | 1.3169587 |
| 40 | $9.72 \times 10^{-4}$ | 0.4420 | 0.44391 | $3.31 \times 10^{-12}$ | 1.3172 | 1.316948 |
| 44 | $4.01 \times 10^{-4}$ |  | 0.44371 | $2.38 \times 10^{-13}$ |  | 1.3176 |

${ }^{a}$ Exact values of $\sigma$ are 0.44357 for $D=1.1$ and 1.3169579 for $D=2$.
to the integral equation, satisfy (4.7) (4.9) with (4.4) and (4.6c'). Here the $d_{q, k}^{(m)}, 0 \leqslant k \leqslant q$, in (4.8) and (4.9) can be determined from (A.5) with $e_{i}(h)=h^{2 i+1}, i=1,2, \ldots$. If we let $n_{l}=2^{l}$, thus $h_{l}=T / 2^{l}, l=0,1,2, \ldots$, as we do in the examples below, then the $d_{q, k}^{(m)}$ are independent of $m$, and can be computed recursively from (A.6) with $\sigma_{i}=2^{-2 i-1}, i=1,2, \ldots$. Thus we find $d_{0,0}^{(m)}=1, \quad d_{1,0}^{(m)}=-1 / 7, \quad d_{1,1}^{(m)}=8 / 7, \quad d_{2,0}^{(m)}=1 / 217, \quad d_{2,1}^{(m)}=-40 / 217, \quad d_{2,2}^{(m)}=$ $256 / 217$, etc. From what has been said in the paragraph following (4.9) and from (A.7), we would expect to have $\Lambda_{q}^{(m)}=O\left(\sigma_{q+1}^{m}\right)=O\left(2^{-(2 q+3) m}\right)$ as $m \rightarrow \infty$, where $\Lambda_{q}^{(m)} \equiv \max _{n_{m} \leqslant j \leqslant n_{m+q}}\left|f\left(x_{j}\right)-\widetilde{f}_{j}\right|$. This is indeed observed numerically.

Example 2. (See Christiansen, 1971).

$$
\int_{0}^{2 \pi} \log \left(2 a \sin \frac{|t-x|}{2}\right) f(x) d x=-\frac{\pi}{2} \cos 2 t
$$

Provided $a \neq 1$, the unique solution to this equation is $f(t)=\cos 2 t$. Otherwise, the solution is $f(t)=\cos 2 t+c, c$ being an arbitrary constant. We observe that the kernel $K(t, x)=\log [2 a \sin (|t-x| / 2)]$ is periodic with period $2 \pi$ in both $t$ and $x$, and is of the form described in ( $c^{\prime}$ ) of Section 4.2, namely, $K(t, x)=H_{1}(t, x) \log |t-x|+\hat{H}_{2}(t, x)$, with $H_{1}(t, t)=1$ and $\hat{H}_{2}(t, t)=\log a$.

Table II gives some of the results obtained with $a=\sqrt{e}$ for $\Lambda_{q}^{(m)}, N=n_{m+q}=2^{m+q}, m+q=3(1) 7$. Note that, for a given row in this
Table II. Results for $\Lambda_{q}^{(m)}$ for Example $2^{a}$

|  | $q$ |  |  |  |  |  |  |  |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| $m+q$ | 0 | 1 | 2 | 3 | 4 | 5 | 6 | 7 |
| 3 | $3.8 \times 10^{-2}$ | $9.9 \times 10^{-3}$ | $4.0 \times 10^{-2}$ | $4.9 \times 10^{-2}$ |  |  |  |  |
| 4 | $4.7 \times 10^{-3}$ | $2.3 \times 10^{-4}$ | $7.4 \times 10^{-5}$ | $3.7 \times 10^{-4}$ | $4.7 \times 10^{-4}$ |  |  |  |
| 5 | $5.9 \times 10^{-4}$ | $6.9 \times 10^{-6}$ | $4.3 \times 10^{-7}$ | $\underline{1.4 \times 10^{-7}}$ | $8.8 \times 10^{-7}$ | $1.1 \times 10^{-6}$ |  |  |
| 6 | $7.3 \times 10^{-5}$ | $2.1 \times 10^{-7}$ | $3.2 \times 10^{-9}$ | $2.1 \times 10^{-10}$ | $6.9 \times 10^{-11}$ | $5.0 \times 10^{-10}$ | $6.3 \times 10^{-10}$ |  |
| 7 | $9.2 \times 10^{-6}$ | $6.6 \times 10^{-9}$ | $2.5 \times 10^{-11}$ | $5.0 \times 10^{-13}$ | $1.5 \times 10^{-13}$ | $\underline{1.2 \times 10^{-13}}$ | $1.8 \times 10^{-13}$ | $1.8 \times 10^{-13}$ |

[^1]table, $N$, the number of abscissas in (4.7)-(4.9), is he same for every member of this row. Thus the first column is the result of no extrapolation, the second, of one extrapolation, the third, of two extrapolations, etc. As can be seen for a given number of abscissas $N=2^{r}$, the best results are obtained roughly for $q \approx r / 2$.

Example 3. (see Henrici, 1979, pp. 492-494). Let $\Gamma: z=z(\tau)$, $0 \leqslant \tau \leqslant \beta$, be the boundary curve of a Jordan region $D$ containing the point $z=0$, and let $f(z)$ be the function mapping $D$ conformally onto the unit disk $|\omega|<1$ in such a manner that $f(0)=0$ and $f^{\prime}(0)>0$. To determine $f(z)$ it is sufficient to know its values on the boundary of $D$. Because $|f(z)|=1$, it suffices to know $\vartheta(\tau)=\arg (f(z(\tau))), 0 \leqslant \tau \leqslant \beta$. Let $\xi(\tau)=\vartheta^{\prime}(\tau)$. Then, provided the capacity of $\Gamma$ is different than $1, \xi(\tau)$ is the unique solution of

$$
\begin{equation*}
\int_{0}^{\beta} \log |z(\sigma)-z(\tau)| \xi(\tau) d \tau=2 \pi \log |z(\sigma)|, \quad 0 \leqslant \sigma \leqslant \beta \tag{5.6}
\end{equation*}
$$

which is known as Symm's equation (see Symm, 1966). The capacity of $\Gamma$ is different from 1 , in particular, if $\Gamma$ is entirely within or entirely without the unit circle.

Obviously, the kernel $K(\sigma, \tau)=\log |z(\sigma)-z(\tau)|$ is periodic in both $\sigma$ and $\tau$ with period $\beta$, and is of the form $K(\sigma, \tau)=H_{1}(\sigma, \tau) \times$ $\log |\sigma-\tau|+\hat{H}_{2}(\sigma, \tau)$ with $H_{1}(\sigma, \sigma)=1$ and $\hat{H}_{2}(\sigma, \sigma)=\log \left|z^{\prime}(\sigma)\right|$. Also, $\log |z(\tau)|$ and $\xi(\tau)$ are periodic with period $\beta$.

Let $\beta=2 \pi$ and let $D$ be the elliptic domain whose boundary curve is $\Gamma: z(\tau)=c\left(e^{i \tau}+\varepsilon e^{-i \tau}\right), 0 \leqslant \tau \leqslant 2 \pi$, with $0<\varepsilon<1$, and $c>0$ chosen so that $\Gamma$ is entirely without the unit circle. (Actually the capacity of $\Gamma$ in this case is $c$, so that it is sufficient to choose $c \neq 1$.) The semiaxes of $D$ are $c(1+\varepsilon)$ and $c(1-\varepsilon)$. The solution for $\xi(\tau)$ can be expressed as

$$
\xi(\tau)=1+4 \sum_{k=1}^{\infty}(-1)^{k} \frac{\varepsilon^{k}}{1+\varepsilon^{2 k}} \cos (2 k \tau)
$$

Note that both $\log |z(\tau)|$ and $\xi(\tau)$ are analytic functions of $\tau$.
Observe that $\xi(\tau)$ for this example is symmetric with respect to both the $\operatorname{Re} z$ and the $\operatorname{Im} z$ axes. This can be utilized to reduce the dimensions of the matrix ( $\tilde{K}_{i j}$ ) by 4 , thus reducing the storage and computing time considerably.

Tables III and IV give some of the results obtained with $c=50$ and $\varepsilon=0.1$ and $\varepsilon=0.5$, respectively, for $\Lambda_{q}^{(m)}, N=n_{m+q}=2^{m+q}, m+q \leqslant 7$. As in Table II, in these tables too, for a given row, $N$, the number of abscissas in (4.7)-(4.9), is the same for every member of this row.

Table III. Results for $\Lambda_{q}^{(m)}$ for Example 3 with $\varepsilon=0.1$ and $c=50^{a}$

| $m+q$ | $q$ |  |  |  |
| :---: | :---: | :---: | :---: | :---: |
|  | 0 | 1 | 2 | 3 |
| 2 | $\underline{1.6 \times 10^{-1}}$ |  |  |  |
| 3 | $2.9 \times 10^{-2}$ | $2.7 \times 10^{-2}$ |  |  |
| 4 | $4.0 \times 10^{-3}$ | $8.1 \times 10^{-4}$ | $4.5 \times 10^{-3}$ |  |
| 5 | $5.0 \times 10^{-4}$ | $2.7 \times 10^{-5}$ | $6.1 \times 10^{-5}$ |  |
| 6 | $6.3 \times 10^{-5}$ | $7.1 \times 10^{-7}$ | $\underline{1.0 \times 10^{-7}}$ | $4.8 \times 10^{-7}$ |
| 7 | $7.8 \times 10^{-6}$ | $2.2 \times 10^{-8}$ | $6.7 \times 10^{-10}$ | $\underline{1.5 \times 10^{-10}}$ |

${ }^{a}$ Here $N=2^{m+q}$ is the number of abscissas in the quadrature method and $q$ is the number of extrapolations in the corresponding numerical quadrature formula. The best result for fixed $N$ (or $m+q$ ) is underlined.

For small values of $\varepsilon$ the ellipse is close to a circle. Therefore, $\xi(\tau)$ does not change rapidly with $\sigma$, and this explains the high accuracy obtained for the approximations to $\xi(\tau)$ even with a small number of abscissas when $\varepsilon=0.1$. For large values of $\varepsilon$, however, the ellipse is elongated, and this leads to rapid changes in $\xi(\tau)$ in the vicinity of $\tau=\pi / 2$ and $\tau=3 \pi / 2$, i.e., where $\xi(\tau)$ is maximal. This then explains the slow convergence of the approximations for $\varepsilon=0.5$. Furthermore, extrapolation becomes effective in this case starting with a relatively large $N$.

To improve the performance of the quadrature method above for the cases in which $\xi(\tau)$ has rapid changes we can make a change of variable of integration $\tau=\tau(\psi)$ so that $\xi(\tau(\psi))$ changes slowly as a function of $\psi$. This

Table IV. Results for $\Lambda_{q}^{(m)}$ for Example 3 with $\varepsilon=0.5$ and $c=50^{a}$

|  |  | $q$ |  |
| :---: | :---: | :---: | :---: |
| $m+q$ | 0 | 1 | 2 |
| 5 | $\frac{5.7 \times 10^{-2}}{1.6 \times 10^{-3}}$ | $1.9 \times 10^{-1}$ |  |
| 6 | $\frac{1.5 \times 10^{-3}}{9.4 \times 10^{-4}}$ | $\underline{3.2 \times 10^{-5}}$ | $3.6 \times 10^{-5}$ |
| 7 |  |  |  |

[^2]can be achieved by picking $\tau(\psi)$ such that $d \tau / d \psi$ becomes small where $d \xi / d \tau$ is large. This is equivalent to having more abscissas in places where $\xi(\tau)$ changes rapidly. Needless to say, the transformation $\tau=\tau(\psi)$ should be such that $d \tau / d \psi$ is a periodic function of $\psi$.

For the example under consideration we can choose

$$
\frac{d \psi}{d \tau}=\frac{\eta\left(1+\eta^{2}\right)^{1 / 2}}{\eta^{2}+\cos ^{2} \tau}, \quad \eta \text { a positive constant }
$$

so that

$$
\psi(\tau)=\tan ^{-1}\left[\frac{\left(1+\eta^{2}\right)^{1 / 2}}{\eta} \tan \tau\right], \quad 0 \leqslant \tau \leqslant \pi / 2
$$

or

$$
\tau(\psi)=\tan ^{-1}\left[\frac{\eta}{\left(1+\eta^{2}\right)^{1 / 2}} \tan \psi\right], \quad 0 \leqslant \psi \leqslant \pi / 2
$$

$\tau(\psi)$ is now extended so that

$$
\begin{array}{ll}
\tau(\psi)=\pi-\tau(\pi-\psi), & \pi / 2 \leqslant \psi \leqslant \pi \\
\tau(\psi)=\pi+\tau(\psi-\pi), & \pi \leqslant \psi \leqslant 2 \pi
\end{array}
$$

We used this transformation with different values of $\eta$. In Table V we give some of the results with $c=50$ and $\varepsilon=0.5$ obtained for the errors at $\tau=\pi / 2$, the point at which the error is maximum, and at $\tau=0$. The number of abscissas in all cases is $N=32$, and $q=0$, i.e., no extrapolation is employed. Nevertheless, the improvement in the results is remarkable.

Table V. Results for the Error at $\tau=0$ and $\tau=\pi / 2$ in the Approximations to $\xi(\tau)$ in Example 3, with $c=50$ and $\varepsilon=0.5$, and the Change of Variable $\tau=\tau(\psi)$ as Described in the Text ${ }^{a}$

| $\eta$ | Error for $\tau=0$ | Error for $\tau=\pi / 2$ |
| :---: | :---: | :---: |
| 0.2 | $4.0 \times 10^{-3}$ | $8.1 \times 10^{-4}$ |
| 0.5 | $1.2 \times 10^{-4}$ | $2.3 \times 10^{-3}$ |
| 100 | $2.4 \times 10^{-4}$ | $5.7 \times 10^{-2}$ |

[^3]
## 6. CONCLUDING REMARKS

In this section we shall briefly discuss some known quadrature methods that are related to those proposed in the present work.

A lot of attention has been paid to Cauchy principal value integrals and singular integral equations. We do not intend to survey all the methods developed for these, but we shall restrict our attention to those that are periodic. When

$$
K(t, x)=\frac{1}{2 \pi} \cot \left(\frac{t-x}{2}\right), \quad 0 \leqslant t, x \leqslant 2 \pi \quad \text { (the Hilbert kernel) }
$$

we have

$$
\widetilde{K}_{i j}=2 h \varepsilon_{i j} K\left(x_{i}, x_{j}\right)=\frac{1}{N} \cot \left[\frac{(i-j) \pi}{2 N}\right] \varepsilon_{i j}, \quad i, j=1, \ldots, 2 N
$$

[cf. (4.1)]. The matrix $\widetilde{K}=\left(\widetilde{K}_{i j}\right)$ is called Wittich's matrix (see Gaier, 1964, p. 76) and its properties are well known. Gutknecht (1981) has used Wittich's matrix to discretize Theodorsen's integral equation for conformal mapping and has analyzed various nonlinear iterative techniques for solving the resulting equations. We note that Theodorsen's equation has the Hilbert kernel as its kernel.

The Hilbert kernel arises as part of the Cauchy kernel on closed curves and Atkinson (1972b) has proposed and analyzed product-type integration formulas for the Cauchy transform

$$
\int_{\Gamma} \frac{\varphi(\zeta)}{\zeta-z} d \zeta, \Gamma \text { closed, } z \in \Gamma
$$

which are different than those proposed in the present work.
As for the weakly singular integral equations, the case that has received the widest attention is that of logarithmic singularity. Periodic integral equations with logarithmically singular kernels arise naturally, for example, in conformal mapping (Symm's equation) and two-dimensional potential theory. Two of the quadrature methods that have been considered for such equations are based on the so-called modified quadrature method (see Kantorovich and Krylov, 1964, p. 102; and Baker, 1977, Chap. 5, Sec. 4). In this method we begin by writing (1.1) in the form

$$
\begin{equation*}
\omega f(t)+\int_{a}^{b} K(t, x)[f(x)-f(t)] d x+f(t) \int_{a}^{b} K(t, x) d x=g(t) \tag{6.1}
\end{equation*}
$$

Employing the numerical quadrature formula $\sum_{j=1}^{n} w_{j} F\left(x_{j}\right)$ to approximate the integral $\int_{a}^{b} F(x) d x$, and using the fact that $\lim _{x \rightarrow t} K(t, x)[f(x)-f(t)]$ $=0$, we replace the integral equation above by the linear equations

$$
\begin{equation*}
\left[\omega+\int_{a}^{b} K\left(x_{i}, x\right) d x\right] \tilde{f}_{i}+\sum_{\substack{j=1 \\ j \neq i}}^{n} w_{j} K\left(x_{i}, x_{j}\right)\left[\tilde{f}_{j}-\tilde{f}_{i}\right]=g\left(x_{i}\right), \quad i=1, \ldots, n \tag{6.2}
\end{equation*}
$$

where $\tilde{f}_{i}$ are approximations to the $f\left(x_{i}\right)$. Kussmaul and Werner (1968) have applied this method with equidistant abscissas $x_{i+1}-x_{i}=h$, $i=1, \ldots, n-1$, to periodic integral equations with logarithmically singular kernels and have shown that the error is $O\left(h^{3}\right)$ as $h \rightarrow 0$, assuming $\int_{a}^{b} K\left(t_{i}, x\right) d x$ have been computed exactly. For kernels of the type $K(p, q)=$ $\log \rho(p, q)$, where $\rho(p, q)=\left\{[x(p)-x(q)]^{2}+[y(p)-y(q)]^{2}\right\}^{1 / 2}$, and $(x(p), y(p)), a \leqslant p \leqslant b$, is the parametric representation of a simple closed curve in the $x-y$ plane, Christiansen (1971) has used the modified quadrature method in (6.1) with $\int_{a}^{b} K\left(t_{i}, q\right) d q$ essentially replaced by a numerical quadrature approximation. The numerical results indicate that this method too has an error of $O\left(h^{3}\right)$ as $h \rightarrow 0$. For both methods above we can show, using Theorem $7 \mathrm{c}^{\prime}$, that the errors in the numerical quadrature formulas are $O\left(h^{3}\right)$ as $h \rightarrow 0$, although this proof for Christiansen's method becomes very complicated.

For weakly singular Fredholm integral equations of the second kind methods based on product integration have also been developed and analyzed by Atkinson (1967, 1972a).

Finally, the approach of this work can easily be extended to coupled Fredholm integral equations for several unknown functions in which several types of singularities occur simultaneously.

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## APPENDIX

In this appendix we summarize some of the important aspects of generalizations of the classical Richardson extrapolation process.

Let the function $A(h)$, where $h>0$ is a continuous or discrete variable, have the asymptotic expansion

$$
\begin{equation*}
A(h) \sim A+\sum_{i=1}^{\infty} \beta_{i} e_{i}(h) \quad \text { as } h \rightarrow 0 \tag{A.1}
\end{equation*}
$$

with $e_{i+1}(h)=o\left(e_{i}(h)\right)$ as $h \rightarrow 0, i=1,2, \ldots A(h)$ and the $e_{i}(h)$ are known, but $A$ and the $\beta_{i}$ are not. We are interested in approximating $A$, which, in general, is $\lim _{h \rightarrow 0} A(h)$ when this limit exists.

Select a sequence of $h$ 's, namely, $h_{0}>h_{1}>\cdots$, such that $\lim _{l \rightarrow \infty} h_{l}=0$, and define $A_{n}^{(j)}$ (the approximation to $A$ ) and the $\bar{\beta}_{i}$ to be the solution of the system of linear equations

$$
\begin{equation*}
A\left(h_{l}\right)=A_{n}^{(j)}+\sum_{i=1}^{n} \bar{\beta}_{i} e_{i}\left(h_{l}\right), \quad j \leqslant l \leqslant j+n \tag{A.2}
\end{equation*}
$$

For some classes of functions $A(h)$ and $e_{i}(h), i=1,2, \ldots$, and some choices of $h_{l}, l=0,1, \ldots$, the following can be shown:
(1) For fixed $n$,

$$
\begin{equation*}
\left.A-A_{n}^{(j)}=O\left(e_{n+1}\left(h_{j}\right)\right) \quad \text { as } j \rightarrow \infty \quad \text { (i.e., } h_{j} \rightarrow 0\right) \tag{A.3}
\end{equation*}
$$

(2) For fixed $j, A_{n}^{(j)} \rightarrow A$ as $n \rightarrow \infty$, and the convergence in this case is better than in the previous case.

For the present work, it is important to note that $A_{n}^{(j)}$ can be expressed as

$$
\begin{equation*}
A_{n}^{(j)}=\sum_{k=0}^{n} d_{n, k}^{(j)} A\left(h_{j+k}\right) \tag{A.4}
\end{equation*}
$$

where the coefficients $d_{n, k}^{(j)}$ can be obtained directly by solving the linear system of equations

$$
\begin{align*}
\sum_{k=0}^{n} d_{n, k}^{(j)} & =1  \tag{A.5}\\
\sum_{k=0}^{n} e_{i}\left(h_{j+k}\right) d_{n, k}^{(j)} & =0, \quad 1 \leqslant i \leqslant n
\end{align*}
$$

The problem of computing the $A_{n}^{(j)}$ recursively has been attacked by several authors. Schneider (1975), Håvie (1979), and Brezinski (1980) have devised an algorithm that has been denoted the $E$ algorithm. Recently, a more efficient algorithm has been derived by Ford and Sidi (1987).

## Special Cases

(a) $e_{i}(h)=h^{\gamma_{i}}, \quad 0<\gamma_{1}<\gamma_{2}<\cdots, \quad \lim _{i \rightarrow \infty} \gamma_{i}=\infty$. For the choice $h_{l}=h_{0} \rho^{l}, l=0,1, \ldots, 0<\rho<1$, the $d_{n, k}^{(j)}$ can be computed from the recursion relation

$$
\begin{equation*}
d_{n, k}^{(j)}=\frac{\sigma_{n} d_{n-1, k}^{(j)}-d_{n-1, k-1}^{(j+1)}}{\sigma_{n}-1}, \quad 0 \leqslant k \leqslant n \tag{A.6}
\end{equation*}
$$

with $d_{n,-1}^{(j)}=0=d_{n, n+1}^{(j)}, n=0,1, \ldots, d_{0,0}^{(j)}=1, j=0,1, \ldots$, and $\sigma_{i}=\rho^{\gamma_{i}}, i=1,2, \ldots$. Obviously, the $d_{n, k}^{(j)}$ are independent of $j$. This development, in slightly different notation, is due to Bulirsch and Stoer (1964), who also give a recursive algorithm for $A_{n}^{(j)}$ and a thorough convergence analysis. First,

$$
\begin{equation*}
A-A_{n}^{(j)}=O\left(\sigma_{n+1}^{j}\right) \quad \text { as } j \rightarrow \infty \tag{A.7}
\end{equation*}
$$

Second, if there exist constants $\widetilde{\beta}_{k}, k=1,2, \ldots$, for which

$$
\begin{equation*}
\left|A(h)-A-\sum_{i=1}^{N-1} \beta_{i} h^{\nu}\right| \leqslant \widetilde{\beta}_{N} h^{N}, \quad h \leqslant h_{0} \tag{A.8}
\end{equation*}
$$

and, for some fixed $R>0$,

$$
\begin{equation*}
\widetilde{\beta}_{k}=O\left(k!R^{k}\right) \quad \text { as } k \rightarrow \infty \tag{A.9}
\end{equation*}
$$

then, for the special case $\gamma_{i}=\gamma_{0}+i r$, for some $\gamma_{0}>0$ and $r>0$,

$$
\begin{equation*}
A-A_{n}^{(j)}=O\left(\omega^{n^{2}}\right) \quad \text { as } n \rightarrow \infty \tag{A.10}
\end{equation*}
$$

where $\omega=\rho^{r / 2}+\varepsilon<1$, any $\varepsilon>0$. For a similar result pertaining to the case of arbitrary $\gamma_{i}$ see Bulirsch and Stoer (1964).
(b) $e_{i}(h)=\varphi(h) h^{i r-r}, r>0$. For arbitrary $h_{l}$, the $A_{n}^{(j)}$ can be computed by using the recursive $W$ algorithm of Sidi (1982). As a consequence of this algorithm we have

$$
\begin{equation*}
d_{n, k}^{(j)}=\frac{\delta_{n, k}^{(j)}}{\sum_{i=0}^{n} \delta_{n, i}^{(j)}}, \quad 0 \leqslant k \leqslant n \tag{A.11}
\end{equation*}
$$

where $\delta_{n, k}^{(j)}$ can be obtained from the recursion relation

$$
\begin{equation*}
\delta_{n, k}^{(j)}=\frac{\delta_{n-1, k-1}^{(j+1)}-\delta_{n-1, k}^{(j)}}{h_{j+n}^{r}-h_{j}^{r}}, \quad 0 \leqslant k \leqslant n \tag{A.12}
\end{equation*}
$$

with $\delta_{n,-1}^{(j)}=0=\delta_{n, n+1}^{(j)}, n=0,1, \ldots$, and $\delta_{0,0}^{(j)}=1 / \varphi\left(h_{j}\right), j=0,1, \ldots$ We note that this is also a special case of the generalized Richardson extrapolation
process of Sidi (1979), a very efficient recursive algorithm for which has recently been given by Ford and Sidi (1987).

We also note that the algorithms given for the two special cases above are more efficient than the algorithms for the general extrapolation algorithms for obtaining $A_{n}^{(j)}$ defined by (A.2), since they take advantage of the special forms of the problem.

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[^1]:    Here $N=2^{m+q}$ is the number of abscissas in the quadrature method and $q$ is the number of extrapolations in the corresponding numerical quadrature formula. The best result for fixed $N$ (or $m+q$ ) is underlined.

[^2]:    ${ }^{a}$ Here $N=2^{m+q}$ is the number of abscissas in the quadrature method and $q$ is the number of extrapolations in the corresponding numerical quadrature formula. The best result for fixed $N$ (or $m+q$ ) is underlined.

[^3]:    ${ }^{a}$ The numerical quadrature formula used has $N=32$ abscissas and does not employ extrapolation. $\xi(0)=$ $0.014671 \ldots$ and $\xi(\pi / 2)=4.5324 \ldots$

