## Applications of Mathematics

## Miroslav Dons; Eva Dontová

A numerical solution of the Dirichlet problem on some special doubly connected regions

Applications of Mathematics, Vol. 43 (1998), No. 1, 53-76

Persistent URL: http://dml.cz/dmlcz/134374

## Terms of use:

© Institute of Mathematics AS CR, 1998

Institute of Mathematics of the Czech Academy of Sciences provides access to digitized documents strictly for personal use. Each copy of any part of this document must contain these Terms of use.


This document has been digitized, optimized for electronic delivery and stamped with digital signature within the project DML-CZ: The Czech Digital Mathematics Library http://dml.cz

# A NUMERICAL SOLUTION OF THE DIRICHLET PROBLEM ON SOME SPECIAL DOUBLY CONNECTED REGIONS 

Miroslav Dont, Eva Dontová, Praha

(Received October 15, 1997)

Abstract. The aim of this paper is to give a convergence proof of a numerical method for the Dirichlet problem on doubly connected plane regions using the method of reflection across the exterior boundary curve (which is analytic) combined with integral equations extended over the interior boundary curve (which may be irregular with infinitely many angular points).

Keywords: Dirichlet problem, integral equations, numerical method
MSC 2000: 31A25, 31A20, 45B05, 65R20, 65N99

## 1. Integral equations

We shall identify the real plane $\mathbb{R}^{2}$ with the complex plane $\mathbb{C}$ in the usual sense. A mapping from a subset of $\mathbb{R}^{2}$ into $\mathbb{R}^{2}$ is then regarded as a complex function of the complex variable.

For the reduction of two integral equations for the Dirichlet problem on doubly connected regions to only one equation we use the so-called reflection function. Note that J. M. Sloss in [20] proved the existence of reflection functions for a sufficiently large family of smooth curves. In [21] he used this reflection function for treating the Dirichlet problem on multiply connected regions bounded by curves with continuously varying curvature. In [4], [5] a suitable integral representation for the solution of the Dirichlet and the Neumann problems on regions with nonsmooth boundaries was established. Let us only recall that if $L$ is an analytic Jordan curve (in $\mathbb{R}^{2}$ ) then,

[^0]under some additional conditions on $L$, there are
$$
R_{g} \text {-an open connected neighbourhood of } L
$$
and
$$
g \text {-a function holomorphic and one-to-one on } R_{g}
$$
such that
\[

$$
\begin{align*}
& \bar{g}\left(R_{g} \cap \operatorname{Int} L\right) \subset \operatorname{Ext} L,  \tag{1.1}\\
& \bar{g}\left(R_{g} \cap \operatorname{Ext} L\right) \subset \operatorname{Int} L,  \tag{1.2}\\
& \bar{g}(z)=z \quad \text { for } z \in L,  \tag{1.3}\\
& \bar{g}(\bar{g}(z))=z \quad \text { for } z \in R_{g} \tag{1.4}
\end{align*}
$$
\]

(see [20] or [4], [5]); here $\bar{g}$ is the complex conjugate of $g$, $\operatorname{Int} L$ and Ext $L$ denote the bounded and the unbounded complementary domain of $L$, respectively. The function $\bar{g}$ is then called the reflection function of the curve $L$. In the sequel we will always suppose that $L$ is an analytic Jordan curve having a reflection function $\bar{g}$ defined on a neighbourhood $R_{g}$ of $L$.

Let us recall the following elementary property of the reflection function (see [4], remark 1.3). Let $h$ be a function harmonic on $R_{g} \cap \operatorname{Int} L$ the first partial derivatives of which are continuously extendible from $R_{g} \cap \operatorname{Int} L$ to $\left(R_{g} \cap \operatorname{Int} L\right) \cup L$; its composition with $\bar{g}$, to be denoted by $h * \bar{g}$, is defined on $R_{g} \cap \operatorname{Ext} L$ and is continuously extendible to $\left(R_{g} \cap \operatorname{Ext} L\right) \cup L$. If $\mathbf{n}_{e}$ denotes the exterior normal of $\operatorname{Int} L$ on $L$ and $\mathbf{n}_{i}$ the corresponding interior normal then (on $L$ )

$$
\begin{equation*}
\frac{\partial h}{\partial \mathbf{n}_{e}}=\frac{\partial(h * \bar{g})}{\partial \mathbf{n}_{i}}=-\frac{\partial(h * \bar{g})}{\partial \mathbf{n}_{e}} . \tag{1.5}
\end{equation*}
$$

Let $K$ be a rectifiable Jordan curve such that

$$
\begin{equation*}
K \subset R_{g} \cap \operatorname{Int} L \tag{1.6}
\end{equation*}
$$

and suppose that

$$
\operatorname{Int} L \cap \operatorname{Ext} K \subset R_{g} .
$$

We shall consider the Dirichlet problem on the domain

$$
\begin{equation*}
S^{+}=\operatorname{Int} L \cap \operatorname{Ext} K \tag{1.7}
\end{equation*}
$$

Further let us denote

$$
\begin{gather*}
S^{-}=\bar{g}\left(S^{+}\right), \quad S=S^{+} \cup L \cup S^{-},  \tag{1.8}\\
\widehat{K}=\bar{g}(K) . \tag{1.9}
\end{gather*}
$$

Then $\widehat{K}$ is a Jordan curve (also rectifiable),

$$
K \cup L \subset \operatorname{Int} \widehat{K}
$$

and further

$$
\begin{gathered}
S^{-}=\operatorname{Ext} L \cap \operatorname{Int} \widehat{K}, \quad S=\operatorname{Ext} K \cap \operatorname{Int} \widehat{K}, \\
\partial S^{+}=K \cup L, \quad \partial S^{-}=L \cup \widehat{K}, \quad \partial S=K \cup \widehat{K} .
\end{gathered}
$$

Suppose that $K$ has a parameterization $\psi$ defined on an interval $\langle a, b\rangle$ [that is, $K=\psi(\langle a, b\rangle)$ and $\psi\left(t_{1}\right) \neq \psi\left(t_{2}\right)$ for any $t_{1}, t_{2} \in\langle a, b\rangle$ with $\left.0<\left|t_{1}-t_{2}\right|<b-a\right]$. For $z \in \mathbb{R}^{2}$ let $\vartheta_{z}$ stand for a single-valued (continuous) branch of the argument of $[\psi-z]$ on the set $\langle a, b\rangle \backslash \psi^{-1}(z)$. For $0<r \leqslant \infty$ let $\gamma_{z, r}$ be the family of all connected components of the set

$$
\{t \in\langle a, b\rangle|0<|\psi(t)-z|<r\}
$$

$\left(\gamma_{z} \equiv \gamma_{z, \infty}\right)$.
The so-called cyclic variation of $\psi$ at $z \in \mathbb{R}^{2}$ is denoted by $v_{r}^{\psi}(z) ; v^{\psi}(z) \equiv v_{\infty}^{\psi}(z)$. Let us recall that

$$
\begin{equation*}
v_{r}^{\psi}(z)=\sum_{I \in \gamma_{z, r}} \operatorname{var}\left[\vartheta_{z} ; I\right] \tag{1.10}
\end{equation*}
$$

where we denote as usual by $\operatorname{var}[\vartheta ; I]$ the total variation of $\vartheta$ on $I$. We will always suppose that

$$
\begin{equation*}
V_{K} \equiv \sup _{z \in K} v^{\psi}(z)<\infty \tag{1.11}
\end{equation*}
$$

Under this assumption $v^{\psi}$ is bounded on $\mathbb{R}^{2}$ and the double layer potential $W_{\psi}(\cdot, f)=$ $W_{K}(\cdot, f)$ for $f \in \mathscr{C}(K)$ is defined on $\mathbb{R}^{2}$ by

$$
\begin{equation*}
W_{K}(z, f)=\frac{1}{\pi} \sum_{I \in \gamma_{z}} \int_{I} f(\psi(t)) \mathrm{d} \vartheta_{z}(t) \tag{1.12}
\end{equation*}
$$

$\left(z \in \mathbb{R}^{2}\right)$. One can similarly define $W_{L}(\cdot, f)$ for $f \in \mathscr{C}(L)$ [an analogue of (1.11) for $L$ is guaranteed since $L$ is analytic]. It is known (see [7], for example) that if (1.11) is fulfilled then finite limits

$$
W_{K}^{e}(\zeta, f) \equiv \lim _{\substack{z \rightarrow \zeta \\ z \in \operatorname{Ext} K}} W_{K}(z, f), \quad W_{K}^{i}(\zeta, f) \equiv \lim _{\substack{z \rightarrow \zeta \\ z \in \operatorname{Int} K}} W_{K}(z, f)
$$

exist for each $\zeta \in K$. Recall how to evaluate these limits. Given $\zeta \in K$ let $t_{1} \in\langle a, b)$, $t_{2} \in(a, b\rangle$ be such that $\zeta=\psi\left(t_{1}\right)=\psi\left(t_{2}\right)$. Due to $v^{\psi}(\zeta)<\infty$ the limits

$$
\begin{align*}
& \tau_{+}(\zeta) \equiv \tau_{+}^{K}(\zeta)=\lim _{t \rightarrow t_{1}+} \frac{\psi(t)-\zeta}{|\psi(t)-\zeta|}  \tag{1.13}\\
& \tau_{-}(\zeta) \equiv \tau_{-}^{K}(\zeta)=\lim _{t \rightarrow t_{2}-} \frac{\psi(t)-\zeta}{|\psi(t)-\zeta|} \tag{1.14}
\end{align*}
$$

exist. Let $\iota$ denote the value of the index of a point from $\operatorname{Int} K$ with respect to $\psi$ (that is, $\iota=1$ if $\psi$ is positively oriented, $\iota=-1$ if $\psi$ is oriented negatively). We say that a vector $\tau \in \mathbb{R}^{2}$ points into $E \subset \mathbb{R}^{2}$ at $\zeta \in \mathbb{R}^{2}$ provided there is a $\delta>0$ such that $\{\zeta+r \tau \mid 0<r<\delta\} \subset E$. To a point $\zeta \in K$ we assign two values $\alpha_{+}(\zeta), \alpha_{-}(\zeta)$ such that

$$
\begin{equation*}
\tau_{+}(\zeta)=\mathrm{e}^{i \alpha_{+}(\zeta)}, \quad \tau_{-}(\zeta)=\mathrm{e}^{i \alpha_{-}(\zeta)} \tag{1.15}
\end{equation*}
$$

and at the same time
(a) $\alpha_{+}(\zeta)<\alpha_{-}(\zeta)<\alpha_{+}(\zeta)+2 \pi$ if $\tau_{+}(\zeta) \neq \tau_{-}(\zeta)$,
(b) $\alpha_{-}(\zeta)=\alpha_{+}(\zeta)+(1-\iota) \pi$ if $\tau_{+}(\zeta)=\tau_{-}(\zeta)$ and the vector $\mathrm{e}^{i\left(\alpha_{+}(\zeta)+\pi\right)}$ points into $\operatorname{Ext} K$ at $\zeta$,
(c) $\alpha_{-}(\zeta)=\alpha_{+}(\zeta)+(1+\iota) \pi$ if $\tau_{+}(\zeta)=\tau_{-}(\zeta)$ and the vector $\mathrm{e}^{i\left(\alpha_{+}(\zeta)+\pi\right)}$ points into $\operatorname{Int} K$ at $\zeta$.

Put

$$
\begin{equation*}
\Delta(\zeta)=\pi-\left(\alpha_{-}(\zeta)-\alpha_{+}(\zeta)\right) \tag{1.16}
\end{equation*}
$$

It is known that, under the condition (1.11),

$$
\begin{align*}
& W_{K}^{e}(\zeta, f)=W_{K}(\zeta, f)-\iota f(\zeta)\left(1-\frac{1}{\pi} \iota \Delta(\zeta)\right)  \tag{1.17}\\
& W_{K}^{i}(\zeta, f)=W_{K}(\zeta, f)+\iota f(\zeta)\left(1+\frac{1}{\pi} \iota \Delta(\zeta)\right) \tag{1.18}
\end{align*}
$$

for $\zeta \in K, f \in \mathscr{C}(K)$, see [8], theorem 2.11, for example. Note that [8] deals with non-tangential limits. For the general case of $\mathbb{R}^{n}, n \geqslant 2$, see [12], [10]. It is easy to
show that the set $\{\zeta \in K||\Delta(\zeta)|>0\}$ formed by the angular points of $K$ is at most countable (cf. 1.7 in [9]).

If we define for $\zeta \in K, f \in \mathscr{C}(K)$

$$
\begin{equation*}
\bar{W}_{K} f(\zeta)=W_{K}(\zeta, f)+\frac{1}{\pi} \Delta(\zeta) f(\zeta) \tag{1.19}
\end{equation*}
$$

then we have

$$
\begin{align*}
W_{K}^{i}(\zeta, f) & =\bar{W}_{K} f(\zeta)+\iota f(\zeta)  \tag{1.20}\\
W_{K}^{e}(\zeta, f) & =\bar{W}_{K} f(\zeta)-\iota f(\zeta)  \tag{1.21}\\
\bar{W}_{K} f(\zeta) & =W_{K}^{i}(\zeta, f)-\iota f(\zeta)=W_{K}^{e}(\zeta, f)+\iota f(\zeta) \tag{1.22}
\end{align*}
$$

Now define

$$
\begin{equation*}
\mathcal{H}_{K} f(z)=W_{K}(z, f)-W_{K}(\bar{g}(z), f) \tag{1.23}
\end{equation*}
$$

for $f \in \mathscr{C}(K), z \in S^{+}$. Since $\bar{g}(\zeta)=\zeta$ for $\zeta \in L$ and $W_{K}(\cdot, f)$ is continuous outside $K$, we have for $\zeta \in L$

$$
\begin{equation*}
\lim _{\substack{z \rightarrow \zeta \\ z \in S^{+}}} \mathcal{H}_{K}=W_{K}(\zeta, f)-W_{K}(\zeta, f)=0 \tag{1.24}
\end{equation*}
$$

For $\zeta \in K$ we obtain

$$
\begin{align*}
\mathcal{H}_{K}^{e} f(z) & \equiv \lim _{\substack{z \rightarrow \zeta \\
z \in S^{+}}} \mathcal{H}_{K} f(z)=W_{K}^{e}(\zeta, f)-W_{K}(\bar{g}(\zeta), f)  \tag{1.25}\\
& =\bar{W}_{K} f(\zeta)-\iota f(\zeta)-W_{K}(\bar{g}(\zeta), f)
\end{align*}
$$

Consider now the equation

$$
\begin{equation*}
\mathcal{H}_{K}^{e} f(\zeta)=u(\zeta), \quad \zeta \in K \tag{1.26}
\end{equation*}
$$

where $u \in \mathscr{C}(K)$ is a given function, $f \in \mathscr{C}(K)$ unknown. If $f$ is a solution of (1.26) (for a given $u$ ) then it follows from (1.25), (1.24) that the function $\mathcal{H}_{K}$ (defined on $S^{+}$) is the solution of the Dirichlet problem on $S^{+}$with the boundary condition $u$ on $K$ and the zero boundary condition on $L$.

Denoting

$$
\begin{equation*}
\overline{\mathcal{H}} f(\zeta) \equiv \bar{W}_{K} f(\zeta)-W_{K}(\bar{g}(\zeta), f) \tag{1.27}
\end{equation*}
$$

for $f \in \mathscr{C}(K), \zeta \in K$, we can write (1.26) in the form

$$
\begin{equation*}
(I-\iota \overline{\mathcal{H}}) f=-\iota u \tag{1.28}
\end{equation*}
$$

[see (1.25), (1.22)]. The problem of solvability of (1.28) was investigated in [4]. It has been shown that if the condition

$$
\begin{equation*}
\frac{1}{\pi} \lim _{r \rightarrow 0+} \sup _{\zeta \in K} v_{r}^{\psi}(\zeta)<1 \tag{1.29}
\end{equation*}
$$

is fulfilled then the space of solutions [in $\mathscr{C}(K)$ ] of the homogeneous equation $(I-\iota \overline{\mathcal{H}}) f=0$ is equal to the space of all functions constant on $K$. In [4] a necessary and sufficient condition on $u \in \mathscr{C}(K)$ guaranteeing solvability of (1.28) [and equivalently of (1.26)] is given (the solution is then determined uniquely up to a constant-see [4], theorem 2.16).

Given $z \in \mathbb{R}^{2}$ put $h_{z}(z)=+\infty$,

$$
\begin{equation*}
h_{z}(\zeta)=\frac{1}{\pi} \ln \frac{1}{|\zeta-z|}, \quad \zeta \in \mathbb{R}^{2} \backslash\{z\} \tag{1.30}
\end{equation*}
$$

Fix $z_{0} \in \operatorname{Int} K$ and put for $z \in R_{g}$

$$
\begin{equation*}
v(z)=h_{z_{0}}(z)-h_{z_{0}}(\bar{g}(z)) \tag{1.31}
\end{equation*}
$$

As we have already noted one can define $W_{L}(\cdot, f)$ for $f \in \mathscr{C}(L)$ analogously to the definition of $W_{K}(\cdot, f)$ for $f \in \mathscr{C}(K)$. Analogously we also define

$$
\begin{equation*}
\mathcal{H}_{L} f(z)=W_{L}(z, f)-W_{L}(\bar{g}(z), f) \tag{1.32}
\end{equation*}
$$

for $f \in \mathscr{C}(L), z \in S^{+}$. We will always suppose that the curve $L$ is positively oriented. Then

$$
\begin{equation*}
\lim _{\substack{z \rightarrow \zeta \\ z \in S^{+}}} \mathcal{H}_{L} f(z)=2 f(\zeta) \tag{1.33}
\end{equation*}
$$

for $f \in \mathscr{C}(L), \zeta \in L$ (see [4], section 2.15).
As shown in [4] the solution of the Dirichlet problem on $S^{+}$with boundary conditions $u_{K} \in \mathscr{C}(K)$ on $K$ and $u_{L} \in \mathscr{C}(L)$ on $L$ can be found in the form

$$
h(z)=\mathcal{H}_{K} f_{K}(z)+\mathcal{H}_{L} f_{L}(z)+a v(z)
$$

where $f_{K} \in \mathscr{C}(K), f_{L} \in \mathscr{C}(L), a \in \mathbb{R}$.

Since $v(\zeta)=0$ for $\zeta \in L$ and since (1.24) is valid, it suffices to choose $f_{L}=\frac{1}{2} u_{L}$ due to (1.33). Then it suffices to choose $f_{K}$ as a solution of the equation

$$
\begin{equation*}
(I-\iota \overline{\mathcal{H}}) f_{K}=-\iota\left(u_{K}-\left.\frac{1}{2} \mathcal{H}_{L} u_{L}\right|_{K}-\left.a \cdot v\right|_{K}\right) \tag{1.34}
\end{equation*}
$$

Here $a \in \mathbb{R}$ has to be chosen in such a way that the equation (1.34) be solvable. It is shown in [4] that such $a \in \mathbb{R}$ exists and that it is determined uniquely. For our purposes we formulate the result in the following form.

Theorem 1.1. Suppose that (1.29) is fulfilled. Then for each boundary conditions $u_{k} \in \mathscr{C}(K), u_{L} \in \mathscr{C}(L)$ there are $a \in \mathbb{R}, f_{K} \in \mathscr{C}(K)$ such that the function $h$ of the form

$$
\begin{equation*}
h(z)=\mathcal{H}_{K} f_{K}(z)+\frac{1}{2} \mathcal{H}_{L} u_{L}(z)+a \cdot v(z) \tag{1.35}
\end{equation*}
$$

solves on $S^{+}$the Dirichlet problem with the boundary conditions $u_{K}, u_{L}$. The constant $a$ is determined uniquely, the function $f_{K}$ is determined uniquely up to a constant; $f_{K}$ is a solution of the equation (1.34).

## 2. Abstract form of a convergence theorem

In this part we shall prove one assertion which can be viewed as a convergence theorem. Using this assertion we shall prove the convergence of a numerical method for solving the equation (1.34). The idea comes from W. L. Wendland [25] (in fact, the assertion is only a variant of Theorem 3.5 from [25]; see also [24]).

Lemma 2.1. Given a normed linear space $L$ let

$$
B_{n}: L \rightarrow L, \quad B: L \rightarrow L \quad(n=1,2,3, \ldots)
$$

be linear operators. Suppose that $B_{n}^{-1}, B^{-1}$ exist, $B_{n}^{-1}$ are bounded and there is an $M \in \mathbb{R}$ such that $\left\|B_{n}^{-1}\right\| \leqslant M$. Let $L_{0}$ be a closed subspace of $L, B\left(L_{0}\right)=L_{0}$. Suppose that

$$
\begin{equation*}
B_{n} x \rightarrow B x \tag{2.1}
\end{equation*}
$$

for each $x \in L_{0}$. If $x_{n} \in L, x_{n} \rightarrow x \in L_{0}$, then

$$
\begin{equation*}
B_{n}^{-1} x_{n} \rightarrow B^{-1} x \tag{2.2}
\end{equation*}
$$

In particular: Let $L$ be a Banach space, $A_{n}, A: L \rightarrow L$ bounded linear operators and suppose that there is a $\lambda \in \mathbb{R}, \lambda<1$, such that

$$
\begin{equation*}
\left\|A_{n}\right\| \leqslant \lambda, \quad\|A\| \leqslant \lambda \tag{2.3}
\end{equation*}
$$

Let $L_{0}$ be a closed subspace of $L,(I-A)\left(L_{0}\right)=L_{0}$. Suppose that

$$
\begin{equation*}
A_{n} x \rightarrow A x \tag{2.4}
\end{equation*}
$$

for each $x \in L_{0}$. Then

$$
\begin{equation*}
\left(I-A_{n}\right)^{-1} x_{n} \rightarrow(I-A)^{-1} x \tag{2.5}
\end{equation*}
$$

whenever $x_{n} \in L, x_{n} \rightarrow x \in L_{0}$.
Proof. We have

$$
B_{n}^{-1}-B^{-1}=B_{n}^{-1}\left(B-B_{n}\right) B^{-1}
$$

Let $x \in L_{0}, \varepsilon>0$. Then $B^{-1} x \in L_{0}\left[\right.$ as $B\left(L_{0}\right)=L_{0}$ by the assumption] and it follows from (2.1) that there is $n_{0}$ such that

$$
\left\|\left(B-B_{n}\right)\left(B^{-1} x\right)\right\|<\frac{\varepsilon}{M}
$$

for $n>n_{0}$. Hence for $n>n_{0}$ we get

$$
\left\|B_{n}^{-1} x-B^{-1}\right\|=\left\|B_{n}^{-1}\left(B-B_{n}\right)\left(B^{-1} x\right)\right\| \leqslant\left\|B_{n}^{-1}\right\|\left\|\left(B-B_{n}\right)\left(B^{-1} x\right)\right\|<\varepsilon
$$

and thus

$$
B_{n}^{-1} x \rightarrow B^{-1} x
$$

If now $x_{n} \in L, x_{n} \rightarrow x \in L_{0}$, then

$$
\begin{aligned}
\left\|B_{n}^{-1} x_{n}-B^{-1} x\right\| & \leqslant\left\|B_{n}^{-1} x_{n}-B_{n}^{-1} x\right\|+\left\|B_{n}^{-1} x-B^{-1} x\right\| \\
& \leqslant M\left\|x_{n}-x\right\|+\left\|B_{n}^{-1} x-B^{-1} x\right\| \rightarrow 0
\end{aligned}
$$

and the first part of Lemma is proved.
The second part follows immediately from the first.

Note that the above assertion is the Lemma 3.1 from [2] slightly modified for our purposes. The following assertion is the Lemma 3.2 from [2].

Lemma 2.2. Let $X$ be a Banach space, $X_{0} \subset X$, a closed subspace. Let $Q, B$ be bounded linear operators,

$$
Q: X \rightarrow X, \quad B: X \rightarrow X_{0}
$$

let $\|Q\|<1$, and suppose that $Q: X_{0} \rightarrow X_{0}$. Then

$$
\begin{equation*}
(I-Q-B)^{-1}(0) \subset X_{0} \tag{2.6}
\end{equation*}
$$

Suppose in addition that $B$ is compact. If for each $f \in X_{0}$ the equation

$$
\begin{equation*}
(I-Q-B) g=f \tag{2.7}
\end{equation*}
$$

(in unknown $g$ ) has a unique solution in $X_{0}$ then (2.7) is uniquely solvable in $X$ for each $f \in X$.

Lemma 2.3. Let $X$ be a Banach space, $X_{0} \subset X$ its closed subspace. Let $Q, B$ : $X \rightarrow X$ be bounded linear operators and suppose that $B$ is compact, $\|Q\|<1$ and

$$
B: X \rightarrow X_{0}, \quad Q: X_{0} \rightarrow X_{0}
$$

Let $H_{n} \subset X(n=1,2,3, \ldots)$ be subspaces in $X$ and let

$$
P_{n}: X \rightarrow H_{n}
$$

be projections, $\left\|P_{n}\right\|=1$, and suppose that for each $f \in X_{0}$

$$
\begin{equation*}
\left\|P_{n} f-f\right\| \rightarrow 0 \quad \text { for } n \rightarrow \infty \tag{2.8}
\end{equation*}
$$

Further let $B_{n}$ be compact operators, $B_{n}: X \rightarrow X_{0}$, and suppose that $B_{n}$ are collectively compact (which means that the set

$$
\begin{equation*}
\left\{B_{n} f \mid n \in \mathbb{N}, f \in X,\|f\| \leqslant 1\right\} \tag{2.9}
\end{equation*}
$$

is relatively compact) and that for each $f \in X_{0}$

$$
\begin{equation*}
B_{n} f \rightarrow B f \quad \text { for } n \rightarrow \infty \tag{2.10}
\end{equation*}
$$

Consider the equations

$$
\begin{align*}
(I-Q-B) u & =f,  \tag{2.11}\\
\left(I-Q P_{n}-B_{n}\right) u_{n} & =f, \tag{2.12}
\end{align*}
$$

where $f \in X$ is given and $u, u_{n} \in X$ are unknown. Suppose that for each $f \in X_{0}$ the equation (2.11) has a unique solution in $X_{0}$. Then there is $n_{0}$ such that for each $n>n_{0}$ and each $f \in X$ the equation (2.12) is uniquely solvable in $X$. At the same time there are constants $c_{1}, c_{2}$ such that the corresponding solutions of (2.11), (2.12) satisfy the estimates

$$
\begin{gather*}
\left\|u_{n}\right\| \leqslant c_{1}\|u\| \leqslant c_{2}\|f\|  \tag{2.13}\\
\left\|u-u_{n}\right\| \leqslant c_{2}\left(\left\|Q P_{n} u-Q u\right\|+\left\|B_{n} u-B u\right\|\right) . \tag{2.14}
\end{gather*}
$$

Proof. First note that it follows from Lemma 2.2 that under the given assumptions the equation (2.11) is uniquely solvable not only in $X_{0}$ but also in $X$ (for $f \in X)$.

Let us show that for all sufficiently large $n$ also the equation (2.12) is uniquely solvable in $X$. Since $B_{n}$ is compact and $\left\|Q P_{n}\right\|<1$ (as $\|Q\|<1,\left\|P_{n}\right\|=1$ ), by the Riesz-Schauder theory it suffices to show that the homogeneous equation

$$
\begin{equation*}
\left(I-Q P_{n}-B_{n}\right) u_{n}=0 \tag{2.15}
\end{equation*}
$$

has (in $X$ ) only the trivial solution. Suppose that there are infinitely many $n$ for which (2.15) has a non-trivial solution. One can suppose for simplicity that for each natural $n$ the equation (2.15) possesses a non-trivial solution. Then there are

$$
u_{n} \in X, \quad\left\|u_{n}\right\|=1
$$

such that (2.15) is valid, which can be written in the form

$$
\left(I-Q P_{n}\right) u_{n}=B_{n} u_{n},
$$

that is

$$
\begin{equation*}
u_{n}=\left(I-Q P_{n}\right)^{-1} B_{n} u_{n} \tag{2.16}
\end{equation*}
$$

[as $\left\|Q P_{n}\right\|<1$ the inverse $\left(I-Q P_{n}\right)^{-1}$ exists]. Since $B_{n}$ are collectively compact and $\left\|u_{n}\right\|=1$, there is a subsequence $\left\{u_{n_{k}}\right\}$ such that the sequence $\left\{B_{n_{k}} u_{n_{k}}\right\}$ is
convergent. For simplicity let us denote this subsequence by the same symbol $\left\{u_{n}\right\}$. We can thus write

$$
\begin{equation*}
B_{n} u_{n}=v_{n} \rightarrow v \in X_{0} \tag{2.17}
\end{equation*}
$$

(as by assumption $B_{n}: X \rightarrow X_{0}$ and $X_{0}$ is closed). Using Lemma 2.1 we now obtain

$$
\left(I-Q P_{n}\right)^{-1} v_{n} \rightarrow(I-Q)^{-1} v \in X_{0}
$$

( $Q: X_{0} \rightarrow X_{0}$ by assumption), whence [see (2.16)]

$$
\begin{equation*}
u_{n} \rightarrow(I-Q)^{-1} v=u \in X_{0} \tag{2.18}
\end{equation*}
$$

We have

$$
\begin{equation*}
\left\|B_{n} u_{n}-B u\right\| \leqslant\left\|B_{n} u_{n}-B_{n} u\right\|+\left\|B_{n} u-B u\right\| . \tag{2.19}
\end{equation*}
$$

Since $u \in X_{0}$ it follows from the assumption [see (2.10)] that

$$
\begin{equation*}
\left\|B_{n} u-B u\right\| \rightarrow 0 . \tag{2.20}
\end{equation*}
$$

Since $B_{n}$ are collectively compact they are also uniformly bounded and thus

$$
\begin{equation*}
\left\|B_{n}\left(u_{n}-u\right)\right\| \rightarrow 0 \tag{2.21}
\end{equation*}
$$

Now it follows from (2.21), (2.20) and (2.19) that

$$
v_{n}=B_{n} u_{n} \rightarrow B u
$$

and thus (as $v_{n} \rightarrow v$ )

$$
\begin{equation*}
B u=v . \tag{2.22}
\end{equation*}
$$

Since we have put $u=(I-Q)^{-1} v$ [see (2.18)] we thus obtain

$$
u=(I-Q)^{-1} B u,
$$

so that

$$
\begin{equation*}
(I-Q-B) u=0 \tag{2.23}
\end{equation*}
$$

Since $\left\|u_{n}\right\|=1$ (and $u_{n} \rightarrow u$ ), we have also $\|u\|=1$, which contradicts the unique solvability of the equation (2.11). The first part of the assertion is proved.

Now we shall prove (2.13). Existence of $c_{2}$ in the second inequality in (2.13) is clear since under the given assumptions the operator $(I-Q-B)$ has a bounded inverse.

Suppose that the first inequality in (2.13) is not valid (that is there exists no $c_{1}$ such that $\ldots$ ). Then there are $f_{n} \in X, u_{n} \in X, g_{n} \in X$ (more precisely, there is a subsequence of indices $n$ ) such that

$$
\begin{align*}
(I-Q-B) u_{n} & =f_{n},  \tag{2.24}\\
\left(I-Q P_{n}-B_{n}\right) g_{n} & =f_{n} \tag{2.25}
\end{align*}
$$

and

$$
\left\|u_{n}\right\|=1, \quad\left\|g_{n}\right\| \rightarrow+\infty
$$

The equalities (2.24), (2.25) yield

$$
\begin{equation*}
(I-Q-B) u_{n}=\left(I-Q P_{n}-B_{n}\right) g_{n} . \tag{2.26}
\end{equation*}
$$

Denoting

$$
h_{n}=\frac{g_{n}}{\left\|g_{n}\right\|}
$$

(2.26) can be written in the form

$$
\left(I-Q P_{n}\right) h_{n}=\frac{1}{\left\|g_{n}\right\|}(I-Q-B) u_{n}+B_{n} h_{n}
$$

or

$$
h_{n}=\left(I-Q P_{n}\right)^{-1}\left\{\frac{1}{\left\|g_{n}\right\|}(I-Q-B) u_{n}+B_{n} h_{n}\right\} .
$$

Since

$$
\left\|u_{n}\right\|=1, \quad\left\|g_{n}\right\| \rightarrow+\infty
$$

we have

$$
\frac{1}{\left\|g_{n}\right\|}(I-Q-B) u_{n} \rightarrow 0
$$

The operators $B_{n}$ are collectively compact by the assumption. As $\left\|h_{n}\right\|=1$ there is a subsequence such that

$$
B_{n} h_{n} \equiv v_{n} \rightarrow v \in X_{0} .
$$

Now it follows from Lemma 2.1 that

$$
h_{n}=\left(I-Q P_{n}\right)^{-1}\left\{\frac{1}{\left\|g_{n}\right\|}(I-Q-B) u_{n}+B_{n} h_{n}\right\} \rightarrow(I-Q)^{-1} v=h \in X_{0} .
$$

In the same way as in the first part of the proof one can find that

$$
B_{n} h_{n} \rightarrow B h
$$

(which means that $v=B h$ ) and obtain that

$$
h=(I-Q)^{-1} B h
$$

and thus

$$
(I-Q-B) h=0 .
$$

But $\|h\|=1\left(\right.$ as $\left.\left\|h_{n}\right\|=1\right), h \in X_{0}$, which contradicts the assumption that (2.11) is uniquely solvable in $X_{0}$.

It suffices to prove (2.14). Consider $n>n_{0}$ [where $n_{0}$ is such that for each $n>n_{0}$ the equation (2.12) is uniquely solvable for any $f \in X]$. Given $f \in X$ let $u$ be the solution of (2.11) [as we have noted, (2.11) is uniquely solvable in $X$ due to Lemma 2.2] and let $u_{n}$ be the solution of (2.12). Then

$$
u-Q u-B u=u_{n}-Q P_{n} u_{n}-B_{n} u_{n},
$$

which can be written in the form

$$
\left(u-u_{n}\right)-Q P_{n}\left(u-u_{n}\right)-B_{n}\left(u-u_{n}\right)=Q u-Q P_{n} u+B u-B_{n} u
$$

that is $\left(u-u_{n}\right)$ is a solution of the equation (2.12) with the right hand side equal to

$$
Q u-Q P_{n} u+B u-B_{n} u .
$$

Now (2.14) follows immediately from the second inequality in (2.13).

Proposition 2.4. Let $X$ be a Banach space, $X_{0} \subset X$ its closed subspace, $q, s \in X_{0}, s \neq 0$. Let $Q, C: X \rightarrow X$ be two bounded linear operators, $C: X \rightarrow X_{0}$, $C$ compact, $Q: X_{0} \rightarrow X_{0},\|Q\|<1$. Let $H_{n}$ be subspaces in $X, P_{n}: X \rightarrow H_{n}$ projections $(n=1,2,3, \ldots)$ such that $\left\|P_{n}\right\|=1, P_{n} s=s(n \in \mathbb{N})$ and suppose that for each $f \in X_{0}$

$$
\left\|P_{n} f-f\right\| \rightarrow 0 \quad \text { for } n \rightarrow \infty
$$

Consider the equations

$$
\begin{align*}
(I-Q-C) u+a q & =f  \tag{2.27}\\
\left(I-Q P_{n}-C P_{n}\right) u_{n}+a_{n} q & =f, \tag{2.28}
\end{align*}
$$

where $f \in X$ is given while $u, u_{n} \in X, a, a_{n} \in \mathbb{R}$ are unknown. Suppose that for each $f \in X_{0}$ the equation (2.27) has a solution ( $\left.u, a\right)$ in $X_{0} \times \mathbb{R}$ where $a \in \mathbb{R}$ is determined uniquely; further suppose that if $\left(u_{0}, a\right)$ is a solution of $(2.27)\left(u_{0} \in X_{0}\right)$ then the set of all solutions of (2.27) is equal to

$$
\begin{equation*}
\left\{\left(u_{0}+t s, a\right) \mid t \in \mathbb{R}\right\} \tag{2.29}
\end{equation*}
$$

Then there is $n_{0}$ such that for each $n>n_{0}$ and for each $f \in X$ the equation (2.28) has a solution in $X \times \mathbb{R}$. The constant $a_{n}$ is then determined uniquely and the set of all solutions of $(2.28)$ is of the form

$$
\begin{equation*}
\left\{\left(u_{n}^{0}+t s, a_{n}\right) \mid t \in \mathbb{R}\right\} \tag{2.30}
\end{equation*}
$$

where $\left(u_{n}^{0}, a_{n}\right)$ is a solution of (2.28).
Further let $r: X \rightarrow \mathbb{R}$ be a bounded linear functional such that $r(s) \neq 0, f \in X_{0}$, $(u, a)$ is a solution of $(2.27)$ such that $r(u)=0$ and let $\left(u_{n}, a_{n}\right)$ be a solution of (2.28) such that $r\left(P_{n} u_{n}\right)=0$. Then

$$
\begin{equation*}
u_{n} \rightarrow u, \quad a_{n} \rightarrow a \tag{2.31}
\end{equation*}
$$

Proof. Let $h$ be an arbitrary bounded linear functional on $X$ such that $h(s) \neq 0$. For $x \in X$ put

$$
B x=C x-h(x) q, \quad B_{n} x=C P_{n} x-h(x) q .
$$

It is easy to see that $B: X \rightarrow X_{0}, B_{n}: X \rightarrow X_{0}, B$ is compact, $B_{n}$ are collectively compact and $B_{n} f \rightarrow B f$ for each $f \in X_{0}$.

For $f \in X_{0}$ consider the equation

$$
\begin{equation*}
(I-Q-B) u=f \tag{2.32}
\end{equation*}
$$

Then $u \in X$ is a solution of (2.32) if and only if $(u, h(u))$ solves (2.27). Since (2.27) is solvable in $X_{0} \times \mathbb{R}$ and the set of all solutions of (2.27) is of the form (2.29), we see that (2.32) is in $X_{0}$ uniquely solvable. Now it follows from Lemma 2.3 that for all sufficiently large $n$ the equation

$$
\begin{equation*}
\left(I-Q P_{n}-B_{n}\right) u_{n}=f \tag{2.33}
\end{equation*}
$$

is uniquely solvable in $X$ (for each $f \in X$ ). If $u_{n}$ is a solution of (2.33) then $\left(u_{n}, h\left(u_{n}\right)\right)$ is a solution of (2.28). The fact that the set of all solutions of (2.27) is of the form (2.29) means that $\{t s \mid t \in \mathbb{R}\}$ is the null space of the operator $(I-Q-C)$. Since $P_{n} s=s$ we see that $s$ is contained also in the null space of $\left(I-Q P_{n}-C P_{n}\right)$.

Suppose that there is an $s_{1} \in X$ such that $s, s_{1}$ are linearly independent and $s_{1}$ is contained in the null space of $\left(I-Q P_{n}-C P_{n}\right)$. Suppose that $h$ was chosen such that, in addition, $h\left(s_{1}\right)=0$. Let $u_{n}$ be a solution of (2.33). Then also $u_{n}+s_{1}$ solves (2.33) since

$$
\begin{aligned}
\left(I-Q P_{n}-B_{n}\right)\left(u_{n}+s_{1}\right) & =\left(I-Q P_{n}-C P_{n}\right)\left(u_{n}+s_{1}\right)+h\left(u_{n}+s_{1}\right) q \\
& =\left(I-Q P_{n}-C P_{n}\right) u_{n}+h\left(u_{n}\right) q=f
\end{aligned}
$$

this contradicts the fact that (2.33) has a unique solution. Now we see that the operator $\left(I-Q P_{n}-C P_{n}\right)$ has the same null space as the operator $(I-Q-P)$ which implies that the set of all solutions of (2.28) is of the form (2.30), indeed.

Given $f \in X_{0}$ let $\widetilde{u}$ be a solution of (2.32), $\widetilde{u}_{n}$ a solution of (2.33). From the inequality (2.14) in Lemma 2.3 we obtain that for some $c_{2} \in \mathbb{R}$

$$
\left\|\widetilde{u}-\widetilde{u}_{n}\right\| \leqslant c_{2}\left(\left\|Q\left(P_{n} \widetilde{u}-\widetilde{u}\right)\right\|+\left\|B_{n} \widetilde{u}-B \widetilde{u}\right\|\right)
$$

Since $\widetilde{u} \in X_{0}$ (due to the assumption $f \in X_{0}$ ) we have

$$
P_{n} \widetilde{u} \rightarrow \widetilde{u}, \quad B_{n} \widetilde{u} \rightarrow B \widetilde{u}
$$

and thus

$$
\widetilde{u}_{n} \rightarrow \widetilde{u} .
$$

If $(u, a)$ is a solution of (2.27) and $\left(u_{n}, a\right)$ is a solution of (2.28) then we have

$$
a_{n}=h\left(\widetilde{u}_{n}\right) \rightarrow h(\widetilde{u})=a .
$$

Let $u, u_{n}$ be, in addition, such that $r(u)=0, r\left(P_{n} u_{n}\right)=0$. By the above argument $u, u_{n}$ are of the form

$$
u=\widetilde{u}+t \cdot s, \quad u_{n}=\widetilde{u}_{n}+t_{n} \cdot s
$$

where

$$
t=-\frac{r(\widetilde{u})}{r(s)}, \quad t_{n}=-\frac{r\left(P_{n} \widetilde{u}\right)}{r(s)}
$$

We have

$$
\left\|P_{n} \widetilde{u}_{n}-\widetilde{u}\right\| \leqslant\left\|P_{n} \widetilde{u}_{n}-P_{n} \widetilde{u}\right\|+\left\|P_{n} \widetilde{u}-\widetilde{u}\right\| \leqslant\left\|P_{n}\right\|\left\|\widetilde{u}_{n}-\widetilde{u}\right\|+\left\|P_{n} \widetilde{u}-\widetilde{u}\right\|,
$$

whence $P_{n} \widetilde{u}_{n} \rightarrow \widetilde{u}$ (as $\left\|P_{n}\right\|=1, \widetilde{u}_{n} \rightarrow \widetilde{u}, P_{n} \widetilde{u} \rightarrow \widetilde{u}$ ). Now we see that $t_{n} \rightarrow t$ and we conclude that $u_{n} \rightarrow u$. The assertion is proved.

## 3. The Dirichlet problem

Now consider again the Dirichlet problem on $S^{+}$as described in the first section. We shall look for the solution in the form (1.35). As we have noted the problem can be reduced to the integral equation (1.34).

For simplicity we will suppose throughout this part that $K$ is positively oriented. Then the equation (1.34) can be written in the form

$$
\begin{equation*}
(I-\overline{\mathcal{H}}) f_{K}+a v=g_{0} \tag{3.1}
\end{equation*}
$$

where

$$
\begin{equation*}
g_{0}=-\left(u_{K}-\left.\frac{1}{2} H_{L} u_{L}\right|_{K}\right) \tag{3.2}
\end{equation*}
$$

(instead of $\left.v\right|_{K}$ we write simply $v$ ). We know that the solution of (3.1) is not uniqueit is true that $a \in \mathbb{R}$ is determined uniquely, but $f_{K}$ is determined up to a constant. If we add a condition that the value of $f_{K}$ at a fixed point in $K$ vanishes (for example) then (3.1) will have a unique solution.

One of the ways how to solve (3.1) numerically is to look for a piecewise constant approximation of $f_{K}$. Dividing $K$ into $n$ arcs on which the approximate solution is constant we get a system of $n$ equations with $(n+1)$ unknowns ( $n$ values of $f_{K}$ and the constant $a$ ). If we choose $f_{K}$ vanishing at a fixed point then we can consider a piecewise constant approximation which is equal to zero on a given arc. Then one column in the system of equations can be removed and we obtain a system of $n$ equations with $n$ unknowns (one of which is $a$ ). Solving this system we get an approximation of $f_{K}$ and at the same time the constant $a$ (an approximation of $a$ ).

In the following we describe this method in more detail, and using Proposition 2.4 we will show that for all sufficiently large $n$ the matrix of the mentioned system of (linear) equations is regular and the method converges [under the assumption (1.29)]. The question concerning the rate of convergence is more complicated and is not investigated here.

In the first section $\overline{\mathcal{H}} f, W_{K}(\cdot, f), \bar{W}_{K} f$ etc. were defined for $f \in \mathscr{C}(K)$. Considering the integral on the right hand side of (1.12) to be the Lebesgue-Stieltjes integral (identifying $\vartheta_{z}$ with the corresponding Lebesgue-Stieltjes measure) $W_{K}(\cdot, f)$ is defined whenever $f$ is a bounded Baire function on $K$ (for example). Let $\mathscr{B}(K)$ stand for the space of all bounded Baire functions on $K$ equipped with the supremum norm. If $W_{K}(\cdot, f)$ is defined for $f \in \mathscr{B}(K)$ then also $\bar{W}_{K} f, \overline{\mathcal{H}} f$ are defined for $f \in \mathscr{B}(K)$ [by equalities (1.19), (1.27)], that is $\bar{W}_{K}, \overline{\mathcal{H}}$ can be regarded as operators on $\mathscr{B}(K)$. It is easily seen that

$$
\bar{W}_{K}, \overline{\mathcal{H}}: \mathscr{B}(K) \rightarrow \mathscr{B}(K)
$$

Denote further

$$
\widehat{W}_{K} f(\zeta)=W_{K}(\bar{g}(\zeta), f)
$$

for $f \in \mathscr{B}(K), \zeta \in K$, so that $\overline{\mathcal{H}}=\bar{W}_{K}-\widehat{W}_{K}$.
Lemma 3.1. Let $\operatorname{var}[\psi ;\langle a, b\rangle]<\infty$. Then the operator $\widehat{W}_{K}: f \mapsto \widehat{W}_{K} f$ is compact as an operator on $\mathscr{B}(K)$ and

$$
\widehat{W}_{K}: \mathscr{B}(K) \rightarrow \mathscr{C}(K)
$$

Proof. This follows from the reasoning described in the proof of Lemma 2.2 in [4]. It only suffices to write everywhere $\mathscr{B}(K)$ instead of $\mathscr{C}(K)$.

Now we want to show that the equation (3.1) corresponds to the case described in Proposition 2.4. Denote

$$
X=\mathscr{B}(K), \quad X_{0}=\mathscr{C}(K)
$$

(then $X$ is a Banach space, $X_{0} \subset X$ a closed subspace). The function $v$ defined in Section 1 (more precisely the restriction $\left.v\right|_{K}$ ) will play the role of $q$ in Proposition 2.4; further let $s \equiv 1$ (on $K$ ).

Define the projections $P_{n}$ considered in Proposition 2.4 in the following way. For $n \in \mathbb{N}$ divide $K$ into $n$ disjoint $\operatorname{arcs} K_{i}^{n}$. For example if $K=\psi(\langle a, b\rangle)$ then put for $i=1,2, \ldots, n$

$$
\begin{align*}
I_{i}^{n} & =\left(a+(i-1) \frac{b-a}{n}, a+i \frac{b-a}{n}\right\rangle  \tag{3.3}\\
K_{i}^{n} & =\psi\left(I_{i}^{n}\right) \tag{3.4}
\end{align*}
$$

and choose $z_{i}^{n} \in K_{i}^{n}$-for example

$$
\begin{equation*}
z_{i}^{n}=\psi\left(a+\left(i-\frac{1}{2}\right) \frac{b-a}{n}\right) . \tag{3.5}
\end{equation*}
$$

Suppose that the diameters of the subarcs $K_{i}^{n}$ satisfy

$$
\max _{1 \leqslant i \leqslant n}\left\{\operatorname{diam} K_{i}^{n}\right\} \rightarrow 0 \quad \text { for } n \rightarrow \infty
$$

(if we choose $K_{i}^{n}$ in the way described above this condition is fulfilled, of course).
For $f \in \mathscr{B}(K)$ define $P_{n} f$ as a function which is constant on each arc $K_{i}^{n}$; for $z \in K_{i}^{n}$ put

$$
P_{n} f(z)=f\left(z_{i}^{n}\right) .
$$

It is easily seen that $P_{n} s=s,\left\|P_{n}\right\|=1$ and for $f \in X_{0}[=\mathscr{C}(K)]$ we have

$$
P_{n} f \rightarrow f \quad \text { for } n \rightarrow \infty .
$$

Now express the operator $\overline{\mathcal{H}}$ as a sum of two operators $Q, C$. Denote by $\mathscr{R}$ the set of all $r>0$ for which there is a circumference $S^{r}$ with radius $r$ such that $\mathscr{H}_{1}\left(S^{r} \cap K\right)>0$ (where $\mathscr{H}_{1}$ is the linear Hausdorff measure). It is shown in [9], lemma 1.10, that $\mathscr{R}$ is at most countable. It is shown also in [9] that for any $r>0, r \notin \mathscr{R}$,

$$
\sup _{\zeta \in K}\left(v_{r}^{\psi}(\zeta)+|\Delta(\zeta)|\right)=\sup _{\zeta \in K} v_{r}^{\psi}(\zeta) .
$$

As (1.29) is supposed to be fulfilled we can choose $r_{0}>0, r_{0} \notin \mathscr{R}$, such that

$$
\begin{equation*}
\frac{1}{\pi} \sup _{\zeta \in K}\left(v_{r_{0}}^{\psi}(\zeta)+|\Delta(\zeta)|\right)<1 \tag{3.6}
\end{equation*}
$$

Now put for $\zeta \in K, f \in \mathscr{B}(K)$,

$$
Q f(\zeta)=\frac{1}{\pi}\left(\sum_{I \in \gamma_{\zeta}, r_{0}} \int_{I} f(\psi(t)) \mathrm{d} \vartheta_{\zeta}(t)+\Delta(\zeta) f(\zeta)\right)
$$

where $\gamma_{\zeta, r_{0}}$ has the meaning described before the formula (1.10) and further

$$
\begin{aligned}
D f(\zeta) & =\bar{W}_{K} f(\zeta)-Q f(\zeta) \\
C f(\zeta) & =\overline{\mathcal{H}} f(\zeta)-Q f(\zeta)=\bar{W}_{K} f(\zeta)-Q f(\zeta)-\widehat{W}_{K} f(\zeta) \\
& =D f(\zeta)-\widehat{W}_{K} f(\zeta)
\end{aligned}
$$

It can be shown analogously to the proofs of Lemmas 1.11 and 1.12 in [9] that $D$ maps the unit ball in $\mathscr{B}(K)$ onto a set of functions equicontinuous and uniformly bounded on $K$. Thus $D: X \rightarrow X_{0}$ and $D$ is a compact operator on $X$ and by Lemma 3.1

$$
C: X \rightarrow X_{0}
$$

and $C$ is a compact operator on $X$. Further we see that

$$
Q: X \rightarrow X, \quad Q: X_{0} \rightarrow X_{0}
$$

(since $\bar{W}_{K}: X \rightarrow X, \bar{W}_{K}: X_{0} \rightarrow X_{0}$ ) and $\|Q\|<1$ due to (3.6).
Now we are in a position to apply Proposition 2.4. Under the given notation the equation (3.1) can be written in the form (2.27); the equation (2.28) has now the form

$$
\begin{equation*}
\left(I-\overline{\mathcal{H}} P_{n}\right) u_{n}+a_{n} v=g_{0} . \tag{3.7}
\end{equation*}
$$

Fix $\widetilde{z} \in K$ and for $f \in X$ put

$$
r(f)=f(\widetilde{z})
$$

Now the assumptions of Proposition 2.4 are fulfilled. Suppose, in addition, that $g_{0} \in X_{0}$ [if $g_{0}$ is of the form (3.2) then it suffices to assume that $u_{K} \in \mathscr{C}(K)$ ]. First it follows from Proposition 2.4 that equations (3.1) and (3.7) (for all sufficiently large $n$ ) are solvable in $X$ (even for $g_{0} \in X$ ). If we add the condition $r\left(f_{K}\right)=0$ to (3.1) and the condition $r\left(P_{n} u_{n}\right)=0$ to (3.7) then the equations (3.1), (3.7) are uniquely solvable. If $\left(f_{K}, a\right)$ is the solution of (3.1) and $\left(u_{n}, a_{n}\right)$ is the solution of (3.7) (fulfilling the given conditions) then

$$
u_{n} \rightarrow f_{K}, \quad a_{n} \rightarrow a .
$$

Now we describe the algorithm to which this method leads. Let $I_{i}^{n}, K_{i}^{n}, z_{i}^{n}$ be of the form (3.3), (3.4), (3.5). Put further

$$
\widetilde{z}=\psi(b) \quad(=\psi(a)) ;
$$

then $\widetilde{z} \in K_{n}^{n}$ (for all $n \in \mathbb{N}$ ). Let $\chi_{i}^{n}$ denote the characteristic function of $K_{i}^{n}$ and put

$$
\begin{equation*}
M_{i j}^{n}=\left[(I-\overline{\mathcal{H}}) \chi_{j}^{n}\right]\left(z_{i}^{n}\right) \tag{3.8}
\end{equation*}
$$

for $i=1,2, \ldots, n, j=1,2, \ldots,(n-1)$ and

$$
\begin{equation*}
M_{i n}^{n}=v\left(z_{i}^{n}\right) \tag{3.9}
\end{equation*}
$$

for $i=1,2, \ldots, n$.
Put further

$$
\begin{equation*}
b_{i}^{n}=g_{0}\left(z_{i}^{n}\right) \tag{3.10}
\end{equation*}
$$

for $i=1,2, \ldots, n$,

$$
\begin{equation*}
x_{i}^{n}=u_{n}\left(z_{i}^{n}\right) \tag{3.11}
\end{equation*}
$$

for $i=1,2, \ldots,(n-1)$ and

$$
\begin{equation*}
x_{n}^{n}=a_{n} . \tag{3.12}
\end{equation*}
$$

Let $u_{n}$ be the solution of the equation (3.7) satisfying the condition $r\left(P_{n} u_{n}\right)=0$, that is $u_{n}\left(z_{n}^{n}\right)=0\left(\right.$ as $\left.\widetilde{z} \in K_{n}^{n}\right)$. Then

$$
\begin{aligned}
P_{n} u_{n} & =\sum_{j=1}^{n-1} u_{n}\left(z_{j}^{n}\right) \chi_{j}^{n}=\sum_{j=1}^{n-1} x_{j}^{n} \chi_{j}^{n} \\
\overline{\mathcal{H}} P_{n} u_{n} & =\overline{\mathcal{H}} \sum_{j=1}^{n-1} x_{j}^{n} \chi_{j}^{n}=\sum_{j=1}^{n-1} x_{j}^{n} \overline{\mathcal{H}} \chi_{j}^{n}
\end{aligned}
$$

and the equality (3.7) can be written in the form

$$
u_{n}-\sum_{j=1}^{n-1} x_{j}^{n} \overline{\mathcal{H}} \chi_{j}^{n}+a_{n} v=g_{0}
$$

In particular we get for $z=z_{i}^{n}$

$$
u_{n}\left(z_{i}^{n}\right)-\sum_{j=1}^{n-1} x_{j}^{n}\left(\overline{\mathcal{H}} \chi_{j}^{n}\right)\left(z_{i}^{n}\right)+a_{n} v\left(z_{i}^{n}\right)=g_{0}\left(z_{i}^{n}\right)
$$

Using the above notation we obtain that

$$
\begin{equation*}
\sum_{j=1}^{n} M_{i j}^{n} x_{j}^{n}=b_{i}^{n} \quad(i=1,2, \ldots, n) \tag{3.13}
\end{equation*}
$$

Note that now we see that the matrix $M=\left(M_{i j}^{n}\right)$ is regular (for all sufficiently large $n$ ) since (3.7) and consequently (3.13) are solvable for any right hand side.

Let us show how to evaluate easily the coefficients $M_{i j}^{n}$ (for $j \neq n$ ). First consider the case $i \neq j$. Then $\chi_{j}^{n}\left(z_{i}^{n}\right)=0$ and thus

$$
M_{i j}^{n}=-\overline{\mathcal{H}} \chi_{j}^{n}\left(z_{i}^{n}\right)=W_{K}\left(\bar{g}\left(z_{i}^{n}\right), \chi_{j}^{n}\right)-W_{K}\left(z_{i}^{n}, \chi_{j}^{n}\right)=\frac{1}{\pi}\left(\int_{I_{j}^{n}} \mathrm{~d} \vartheta_{\bar{g}\left(z_{i}^{n}\right)}-\int_{I_{j}^{n}} \mathrm{~d} \vartheta_{z_{i}^{n}}\right) .
$$

The integrals on the right hand side are nothing else but the increments of the argument along the arc $K_{j}^{n}$ with respect to the point $\bar{g}\left(z_{i}^{n}\right)$ (resp. $z_{i}^{n}$ ) which are easy to evaluate [they are equal to the angles "under which the endpoints of $K_{j}^{n}$ are seen from the point $\bar{g}\left(z_{i}^{n}\right)$ or $z_{i}^{n "] . ~}$

Now consider the case $j=i$. Then $z_{i}^{n} \in K_{i}^{n}$ and we get [using equalities (1.19), (1.16) and the fact that $\chi_{i}^{n}\left(z_{i}^{n}\right)=1$ ]

$$
\begin{aligned}
M_{i i}^{n} & =W_{K}\left(\bar{g}\left(z_{i}^{n}\right), \chi_{i}^{n}\right)+1-\bar{W}_{K} \chi_{i}^{n}\left(z_{i}^{n}\right) \\
& =W_{K}\left(\bar{g}\left(z_{i}^{n}\right), \chi_{i}^{n}\right)+1-W_{K}\left(z_{i}^{n}, \chi_{i}^{n}\right)-\frac{1}{\pi} \Delta\left(z_{i}^{n}\right) \\
& =W_{K}\left(\bar{g}\left(z_{i}^{n}\right), \chi_{i}^{n}\right)+1-W_{K}\left(z_{i}^{n}, \chi_{i}^{n}\right)-\frac{1}{\pi}\left[\pi-\left(\alpha_{-}\left(z_{i}^{n}\right)-\alpha_{+}\left(z_{i}^{n}\right)\right)\right] \\
& =W_{K}\left(\bar{g}\left(z_{i}^{n}\right), \chi_{i}^{n}\right)-\left[W_{K}\left(z_{i}^{n}, \chi_{i}^{n}\right)-\frac{1}{\pi} \alpha_{-}\left(z_{i}^{n}\right)+\frac{1}{\pi} \alpha_{+}\left(z_{i}^{n}\right)\right] .
\end{aligned}
$$

Here

$$
W_{K}\left(\bar{g}\left(z_{i}^{n}\right), \chi_{i}^{n}\right)=\frac{1}{\pi} \int_{I_{i}^{n}} \mathrm{~d} \vartheta_{\bar{g}\left(z_{i}^{n}\right)}
$$

For evaluating $W_{K}\left(z_{i}^{n}, \chi_{i}^{n}\right)$ one has to realize that $z_{i}^{n} \in K_{i}^{n}$. Denote for a while

$$
\begin{aligned}
t_{i} & =a+\left(i-\frac{1}{2}\right) \frac{b-a}{n} \\
\mathbf{1}_{I_{i}^{n}} & =\left(a+(i-1) \frac{b-a}{n}, t_{i}\right), \\
\mathbf{2}_{I_{i}^{n}} & =\left(t_{i}, a+i \frac{b-a}{n}\right)
\end{aligned}
$$

Recalling (1.15) we observe that

$$
\lim _{t \rightarrow t_{i}+} \vartheta_{z_{i}^{n}}(t)
$$

coincides with $\alpha_{+}\left(z_{i}^{n}\right)$ up to a multiple of $2 \pi$ and analogously for

$$
\lim _{t \rightarrow t_{i}-} \vartheta_{z_{i}^{n}}(t)
$$

Suppose that $\vartheta_{z_{i}^{n}}$ has been chosen in such a way that

$$
\lim _{t \rightarrow t_{i}+} \vartheta_{z_{i}^{n}}(t)=\alpha_{+}\left(z_{i}^{n}\right), \quad \lim _{t \rightarrow t_{i}-} \vartheta_{z_{i}^{n}}(t)=\alpha_{-}\left(z_{i}^{n}\right) .
$$

Then we have

$$
\begin{aligned}
& W_{K}\left(z_{i}^{n}, \chi_{i}^{n}\right)-\frac{1}{\pi} \alpha_{-}\left(z_{i}^{n}\right)+\frac{1}{\pi} \alpha_{+}\left(z_{i}^{n}\right) \\
&= \frac{1}{\pi}\left\{\int_{\mathbf{1}_{I_{i}^{n}}} \mathrm{~d} \vartheta_{z_{i}^{n}}+\int_{\mathbf{2}_{I_{i}^{n}}} \mathrm{~d} \vartheta_{z_{i}^{n}}-\alpha_{-}\left(z_{i}^{n}\right)+\alpha_{+}\left(z_{i}^{n}\right)\right\} \\
&= \frac{1}{\pi}\left\{\alpha_{-}\left(z_{i}^{n}\right)-\vartheta_{z_{i}^{n}}\left(a+(i-1) \frac{b-a}{n}\right)+\vartheta_{z_{i}^{n}}\left(a+i \frac{b-a}{n}\right)\right. \\
&\left.-\alpha_{+}\left(z_{i}^{n}\right)-\alpha_{-}\left(z_{i}^{n}\right)+\alpha_{+}\left(z_{i}^{n}\right)\right\} \\
&= \vartheta_{z_{i}^{n}}\left(a+i \frac{b-a}{n}\right)-\vartheta_{z_{i}^{n}}\left(a+(i-1) \frac{b-a}{n}\right) .
\end{aligned}
$$

We see that also $M_{i i}^{n}$ can be evaluated easily (analogously to $M_{i j}^{n}$ for $i \neq j$ ). Thus the described numerical method is easy to implement.

Remark. Consider the case $u_{L} \equiv 0$ on $L$ (then $g_{0}=-u_{K}$ ), let $n$ be sufficiently large and let $\left[x_{1}^{n}, x_{2}^{n}, \ldots, x_{n}^{n}\right]$ be the solution of (3.13),

$$
f_{K}^{n}=\sum_{j=1}^{n-1} x_{j} \chi_{j}^{n}
$$

For $z \in S^{+} \cup L$ put

$$
h(z)=\mathcal{H}_{K} f_{K}^{n}(z)+x_{n}^{n} v(z)=\sum_{j=1}^{n-1} x_{j}^{n} \mathcal{H}_{K} \chi_{j}^{n}(z)+x_{n}^{n} v(z)
$$

[the evaluation of $\mathcal{H}_{K} \chi_{j}^{n}(z)$ is quite similar to the evaluation of $M_{i j}^{n}$ for $i \neq j$ ]. Function $h$ is an approximation of the solution of the Dirichlet problem on $S^{+}$with the boundary condition $u_{K}$ on $K$ and $u_{L} \equiv 0$ on $L$. It is easy to see that $h$ has the following properties:
(a) $h$ is harmonic on $S^{+}$,
(b) $h$ is continuos on $S^{+} \cup L$ and $h(z)=0$ for any $z \in L$,
(c) for each $\zeta \in K$ except the endpoints of $K_{i}^{n}(i=1,2, \ldots, n)$ the limit

$$
\lim _{\substack{z \rightarrow \zeta \\ z \in S^{+}}} h(z)
$$

exists and is finite,
(d) for $i=1,2, \ldots, n$

$$
\lim _{\substack{z \rightarrow z_{i}^{n} \\ z \in S^{+}}} h(z)=u_{K}\left(z_{i}^{n}\right)
$$

[1] Ju. D. Burago, V. G. Maz'ja: Some problems of potential theory and theory of functions for domains with nonregular boundaries. Zapisky Naučnych Seminarov LOMI 3 (1967). (In Russian.)
[2] M. Dont: Fourier problem with bounded Baire data. Math. Bohemica 122 (1997), 405-441.
[3] E. Dontová: Reflection and the Dirichlet and Neumann problems. Thesis, Prague, 1990. (In Czech.)
[4] E. Dontová: Reflection and the Dirichlet problem on doubly connected regions. Casopis pro pěst. mat. 113 (1988), 122-147.
[5] E. Dontová: Reflection and the Neumann problem in doubly connected regions. Časopis pro pěst. mat. 113 (1988), 148-168.
[6] Král J.: Some inequalities concerning the cyclic and radial variations of a plane path-curve. Czechoslovak Math. J. 14 (89) (1964), 271-280.
[7] Král J.: On the logarithmic potential of the double distribution. Czechoslovak Math. J. 14 (89) (1964), 306-321.
[8] Král J.: Non-tangential limits of the logarithmic potential. Czechoslovak Math. J. 14 (89) (1964), 455-482.
[9] J. Král: The Fredholm radius of an operator in potential theory. Czechoslovak Math. J. 15 (1965), 454-473, 565-588.
[10] J. Král: Integral Operators in Potential Theory. Lecture Notes in Math. vol. 823, Springer-Verlag 1980.
[11] J. Král: Boundary regularity and normal derivatives of logarithmic potentials. Proc. of the Royal Soc. of Edinburgh 106A (1987), 241-258.
[12] J. Král: The Fredholm method in potential theory. Trans. Amer. Math. Soc. 125 (1966), 511-547.
[13] J. Král: Potential Theory I. State Pedagogic Publishing House, Praha 1965 (in Czech).
[14] J. Král, I. Netuka, J. Veselý: Potential Theory II. State Pedagogic Publishing House, Praha 1972 (in Czech).
[15] V. G. Maz'ja: Boundary Integral Equations, Analysis IV, Encyclopaedia of Mathematical Sciences vol. 27. Springer-Verlag, 1991.
[16] I. Netuka: Double layer potentials and the Dirichlet problem. Czech. Math. J. 24 (1974), 59-73.
[17] I. Netuka: Generalized Robin problem in potential theory. Czech. Math. J. 22 (1972), 312-324.
[18] I. Netuka: An operator connected with the third boundary value problem in potential theory. Czech. Math. J. 22 (1972), 462-489.
[19] I. Netuka: The third boundary value problem in potential theory. Czech. Math. J. 22 (1972), 554-580.
[20] J. M. Sloss: Global reflection for a class of simple closed curves. Pacific J. Math. 52 (1974), 247-260.
[21] J. M. Sloss: The plane Dirichlet problem for certain multiply connected regions. J. Analyse Math. 28 (1975), 86-100.
[22] J. M. Sloss: A new integral equation for certain plane Dirichlet problems. SIAM J. Math. Anal. 6 (1975), 998-1006.
[23] J. M. Sloss, J. C. Bruch: Harmonic approximation with Dirichlet data on doubly connected regions. SIAM J. Numer. Anal. 14 (1977), 994-1005.
[24] F. Stummel: Diskrete Konvergenz linearen Operatoren I and II. Math. Zeitschr. 120 (1971), 231-264.
[25] W. L. Wendland: Boundary element methods and their asymptotic convergence. Lecture Notes of the CISM Summer-School on "Theoretical acoustic and numerical techniques", Int. Centre Mech. Sci., Udine (Italy) (P. Filippi, ed.). Springer-Verlag, Wien, New York, 1983, pp. 137-216.

Authors' addresses: Miroslav Dont, Dept. of Math., Fac. of Electrical Engineering, CTU, Technická 2, 16627 Praha 6, Czech Republic, e-mail dont@math.feld.cvut.cz; Eva Dontová, Dept. of Math., Fac. of Nucl. Sci. \& Phys. Eng., CTU, Břehová 7, 115 19, Praha 1, Czech Republic.


[^0]:    Support of the grant $\mathrm{N}^{\mathrm{o}} 201 / 96 / 0431$ of the Czech Grant Agency is gratefully acknowledged.

