# NONHOMOGENEOUS NAVIER-STOKES PROBLEMS IN $L_{p}$ SOBOLEV SPACES OVER EXTERIOR AND INTERIOR DOMAINS 

G. GRUBB<br>Copenhagen University Mathematics Department, Universitetsparken 5, DK-2100 Copenhagen, Denmark E-mail: grubb@math.ku.dk


#### Abstract

We here present our work on the solvability of completely nonhomogeneous initialboundary value problems for the Navier-Stokes equations, in general anisotropic $L_{p}$ Sobolev and Besov spaces with $p>1$. Introducing a new twist of the method (simplifying slightly), we can now extend the results to exterior domains, for finite time intervals.


## 1 Introduction

In a series of papers, the author has treated the nonhomogeneous Navier-Stokes problem

$$
\begin{align*}
\partial_{t} u-\Delta u+\sum_{j=1}^{n} u_{j} \partial_{j} u+\operatorname{grad} q & =f & & \text { on } Q_{I_{b}}=\Omega \times I_{b}, \\
\operatorname{div} u & =0 & & \text { on } Q_{I_{b}},  \tag{1.1}\\
T_{k}\{u, q\} & =\varphi_{k} & & \text { on } S_{I_{b}}=\Gamma \times I_{b}, \\
r_{0} u & =u_{0} & & \text { on } \Omega
\end{align*}
$$

for bounded domains $\left.\Omega \subset \mathbf{R}^{n}, I_{b}=\right] 0, b\left[\subset \mathbf{R}_{+}\right.$, with various boundary operators $T_{k}$ of Dirichlet, Neumann or intermediate type ( $r_{0}$ indicates restriction to $t=0$; further details are given below in Section 2). Strong solvability results were obtained in anisotropic $L_{2}$ Sobolev spaces in joint works with V. A. Solonnikov $\left[{ }^{11}\right]-\left[{ }^{14}\right]$, and the results have been extended more recently to $L_{p}$ Sobolev spaces $\left[{ }^{6}\right]-\left[{ }^{7}\right]$, that we report on below (in Section 2). Besides this, we give generalizations to exterior domains (in Section 4), based on a simplified proof (in Section 3).

The main technique is to reduce the linearized problem

$$
\begin{align*}
\partial_{t} u-\Delta u+\operatorname{grad} q & =f & & \text { on } Q_{I_{b}}, \\
\operatorname{div} u & =0 & & \text { on } Q_{I_{b}},  \tag{1.2}\\
T_{k}\{u, q\} & =\varphi_{k} & & \text { on } S_{I_{b}}, \\
r_{0} u & =u_{0} & & \text { on } \Omega ;
\end{align*}
$$

which is degenerate parabolic, to a truly parabolic pseudodifferential problem

$$
\begin{align*}
\partial_{t} u-\Delta u+G_{k} u & =f_{k} & & \text { on } Q_{I_{b}}, \\
T_{k}^{\prime} u & =\psi_{k} & & \text { on } S_{I_{b}},  \tag{1.3}\\
r_{0} u & =u_{0} & & \text { on } \Omega ;
\end{align*}
$$

where the general theory of $\left[{ }^{4}\right],\left[{ }^{12}\right],\left[{ }^{9}\right],\left[{ }^{6}\right]$ can be brought into use.
Parabolic problems of the form $\partial_{t} u+A\left(x, D_{x}\right) u=f$ (with initial and boundary conditions) are much harder when $A$ is of pseudodifferential type than when it is a differential operator, since the singularity of the symbol of $A$ at $\xi=0$ has an important effect when there is an extra parameter-dependence (caused by $\partial_{t}$ ). While trying to extend our results to exterior domains, we were inspired by a recent collaboration with R. Seeley [ ${ }^{10}$ ] to look for simplifications in the treatment of (1.3) such that one can take advantage of the fact that the non-differential aspects are connected with the boundary only.

We shall show below in Section 3 how an important step in the treatment of (1.3) can be broken up into three parts, treating: (i) a classical Dirichlet or Neumann heat problem, (ii) a parameter-dependent ps.d.o. problem on the boundary $\Gamma$, (iii) a classical Dirichlet or Neumann problem for the Laplace operator. For exterior problems, this viewpoint has the advantage that we can lean on known results for the unbounded domain, and need the technical ps.d.o. considerations only on the compact manifold $\Gamma$. It gives rather easily some extensions of the results of $\left[{ }^{7}\right]$ to unbounded domains, however for bounded time intervals only.

For the unbounded time interval $\mathbf{R}_{+}$, the results for the Dirichlet problem in [ ${ }^{[7]}$ do not seem readily extendible; and the new method is perhaps too rough. In fact, one may have to work in other spaces than those that we deal with here (e.g. homogeneous spaces or weighted spaces), to get really satisfactory results.

## 2 Results for the interior case

Consider the problems (1.1) and (1.2). Here $u(x, t)$ is the velocity vector $u=\left\{u_{1}, \ldots, u_{n}\right\}, q(x, t)$ is the (scalar) pressure, and $T_{k}$ is one of the following trace operators:

$$
\begin{array}{ll}
T_{0}\{u, q\}=\gamma_{0} u, \\
T_{1}\{u, q\}=\chi_{1} u-\gamma_{0} q \vec{n}, & T_{2}\{u, q\}=\left(\chi_{1} u\right)_{\tau}+\gamma_{0} u_{\nu} \vec{n},  \tag{2.1}\\
T_{3}\{u, q\}=\gamma_{1} u-\gamma_{0} q \vec{n}, & T_{4}\{u, q\}=\gamma_{1} u_{\tau}+\gamma_{0} u_{\nu} \vec{n},
\end{array}
$$

where $\vec{n}=\left(n_{1}, \ldots, n_{n}\right)$ is the (interior) normal at $\Gamma, v_{\nu}$ resp. $v_{\tau}$ denotes the normal resp. tangential component of an $n$-vector field $v$ defined near $\Gamma$ :

$$
\begin{equation*}
v_{\nu}=\vec{n} \cdot v, \quad v_{\tau}=v-(\vec{n} \cdot v) \vec{n}, \tag{2.2}
\end{equation*}
$$

$\gamma_{k} u=\left.\partial_{\nu}^{k} u\right|_{\Gamma}$ (with $\partial_{\nu}=\sum_{j=1}^{n} n_{j} \partial_{j}$ ), and $\chi_{1}$ is the special first order boundary operator defined via the strain tensor $S(u)=\left(\partial_{i} u_{j}+\partial_{j} u_{i}\right)_{i, j=1, \ldots, n}$ as

$$
\begin{equation*}
\chi_{1} u=\gamma_{0} S(u) \vec{n}=\gamma_{0}\left(\sum_{j}\left(\partial_{i} u_{j}+\partial_{j} u_{i}\right) n_{j}\right)_{i=1, \ldots, n} \tag{2.3}
\end{equation*}
$$

For $k=0$ this gives the Dirichlet problem, $k=1$ and 3 give Neumann problems, and $k=2$ and 4 give problems with partially a Dirichlet, partially a Neumann condition. More comments on these boundary conditions in [ ${ }^{14}$ ].

The data are assumed to satisfy

$$
\begin{align*}
& \operatorname{div} u_{0}=0, \quad \text { when } k=1 \text { or } 3 \\
& \operatorname{div} u_{0}=0, \quad \gamma_{0} u_{0, \nu}=0, \varphi_{k, \nu}=0, \quad \text { when } k=0,2 \text { or } 4 \tag{2.4}
\end{align*}
$$

The problem is considered in anisotropic Bessel-potential spaces $H_{p}^{(s, s / 2)}\left(\bar{Q}_{I}\right)^{n}$ and Besov spaces $B_{p}^{(s, s / 2)}\left(\bar{Q}_{I}\right)^{n}$, where, as we recall, the $H_{p}^{(s, s / 2)}$ spaces are generalizations of the integer case

$$
\begin{equation*}
H_{p}^{2 m, m}\left(\bar{Q}_{I_{b}}\right)=\left\{u(x, t) \in L_{p}\left(Q_{I_{b}}\right) \mid D_{x}^{\alpha} D_{t}^{j} u \in L_{p}\left(Q_{I_{b}}\right) \text { for }|\alpha|+2 j \leq 2 m\right\} \tag{2.5}
\end{equation*}
$$

defined via local coordinates and restriction from

$$
\begin{equation*}
H_{p}^{(s, s / 2)}\left(\mathbf{R}^{n} \times \mathbf{R}\right)=\mathrm{OP}\left(\left(|\xi|^{4}+\tau^{2}+1\right)^{-s / 4}\right) L_{p}\left(\mathbf{R}^{n} \times \mathbf{R}\right) \tag{2.6}
\end{equation*}
$$

this scale is preserved under complex interpolation. The Besov scale $B_{p}^{(s, s / 2)}$ is defined slightly differently, but arises from the $H_{p}^{(s, s / 2)}$ scale by suitable real interpolation. (Further details are given e.g. in [ ${ }^{6}$ ].)

The $B_{p}^{(s, s / 2)}$ spaces must be included even if one is mainly interested in solving the problem in spaces (2.6), because they are the correct boundary value spaces, as $\gamma_{j}$ maps $H_{p}^{(s, s / 2)}\left(\bar{Q}_{I_{b}}\right)$ continuously onto the space $B_{p}^{\left(s-j-\frac{1}{p},\left(s-j-\frac{1}{p}\right) / 2\right)}\left(\bar{S}_{I_{b}}\right)$, for $j<s-\frac{1}{p}$. We denote by $\mathcal{B}_{p}^{s+2,(k)}$ the range space for $T_{k}$ applied to $H_{p}^{(s+2, s / 2+1)}\left(\bar{Q}_{I_{b}}\right)^{n}$.

Let us first present the main results of $\left[{ }^{7}\right]$ for bounded domains:
Consider systems of functions

$$
\begin{equation*}
\Phi_{k}=\left\{f, \varphi_{k}, u_{0}\right\} \in H_{p}^{(s, s / 2)}\left(\bar{Q}_{I_{b}}\right)^{n} \times \mathcal{B}_{p}^{s+2,(k)} \times B_{p}^{s+2-2 / p}(\bar{\Omega})^{n} \tag{2.7}
\end{equation*}
$$

for $s>\frac{1}{p}-1$ with $s \geq \frac{n+2}{p}-3$. The system is said to satisfy the compatibility condition of order $s$, when

$$
\begin{array}{ll}
r_{0} \partial_{t}^{l} \varphi_{k, \tau}=\gamma_{0} u_{\tau}^{(l)} & \text { for } k=0,2 l \leq s+2-\frac{3}{p}, \\
r_{0} \partial_{t}^{l} \varphi_{k, \tau}=\left(\chi_{1} u^{(l)}\right)_{\tau} & \text { for } k=1 \text { and } 2,2 l \leq s+1-\frac{3}{p}, \\
r_{0} \partial_{t}^{l} \varphi_{k, \tau}=\gamma_{1} u_{\tau}^{(l)} & \text { for } k=3 \text { and } 4,2 l \leq s+1-\frac{3}{p},
\end{array}
$$

understood as

$$
\begin{array}{ll}
\mathcal{I}\left[\partial_{t}^{l} \varphi_{k, \tau}, u_{\tau}^{(l)}\right]<\infty & \text { if } k=0,2 l=s+2-\frac{3}{p}  \tag{2.8}\\
\mathcal{I}\left[\partial_{t}^{l} \varphi_{k, \tau},\left(S\left(u^{(l)}\right) \vec{n}\right)_{\tau}\right]<\infty & \text { if } k=1 \text { and } 2,2 l=s+1-\frac{3}{p} \\
\mathcal{I}\left[\partial_{t}^{l} \varphi_{k, \tau}, \partial_{\nu} u_{\tau}^{(l)}\right]<\infty & \text { if } k=3 \text { and } 4,2 l=s+1-\frac{3}{p}
\end{array}
$$

here the $u^{(l)}$ are defined successively by

$$
\begin{align*}
u^{(0)} & =u_{0} \\
u^{(l+1)} & =\left(\Delta-G_{k}\right) u^{(l)}-\kappa \sum_{m=0}^{l}\binom{l}{m} \mathcal{Q}_{k}\left(u^{(m)}, u^{(l-m)}\right)+r_{0} \partial_{t}^{l} f_{k} \tag{2.9}
\end{align*}
$$

where the $G_{k}$ are certain singular Green operators stemming from the elimination of the pressure $q$, and

$$
\begin{equation*}
\mathcal{I}[\psi, v]=\int_{t \in I} \int_{x^{\prime} \in \Gamma} \int_{y \in \Omega} \frac{\left|\psi\left(x^{\prime}, t\right)-v(y)\right|^{p}}{\left(\left|x^{\prime}-y\right|^{d}+t\right)^{1+n / d}} d y d \sigma_{x^{\prime}} d t \tag{2.10}
\end{equation*}
$$

We then define the data norm of $\Phi_{k}$ by

$$
\begin{equation*}
\mathcal{N}_{s, p, b}^{(k)}\left(\Phi_{k}\right)=\left(\|f\|_{H_{p}^{(s, s / 2)}\left(\bar{Q}_{I_{b}}\right)^{n}}^{p}+\left\|\varphi_{k}\right\|_{\mathcal{B}_{p}^{s+2,(k)}}^{p}+\left\|u_{0}\right\|_{B_{p}^{s+2-2 / p}(\bar{\Omega})^{n}}^{p}+\mathcal{I}_{s, p, b}\right)^{\frac{1}{p}} \tag{2.11}
\end{equation*}
$$

where $\mathcal{I}_{s, p, b}=0$ if $s+2-\frac{3}{p} \notin \mathbf{N}$, and otherwise equals the possible $\mathcal{I}$ term entering in the compatibility condition. The following result on uniqueness and on the existence of solutions on large time-intervals for small enough data, and on small enough time-intervals for large data, is proved in detail in $\left[{ }^{7}\right]$.

Theorem 2.1 Let $\Omega$ be a smooth bounded open set in $\mathbf{R}^{n}$. Let $k=0,1,2,3$ or 4 , let $s>\frac{1}{p}-1$ with $s \geq \frac{n+2}{p}-3$, and let $b \in \mathbf{R}_{+}$. Consider $\Phi_{k}$ as in (7), satisfying the compatibility condition of order $s$.
$1^{\circ}$ There is at most one solution $\{u, q\}$ with

$$
\begin{equation*}
\{u, \operatorname{grad} q\} \in H_{p}^{(s+2, s / 2+1)}\left(\bar{Q}_{I_{b}}\right)^{n} \times H_{p}^{(s, s / 2)}\left(\bar{Q}_{I_{b}}\right) \tag{2.12}
\end{equation*}
$$

of the Navier-Stokes problem (1.1) for each set of data $\Phi_{k}$ (where $q$ for $k=0,2$ or 4 is subject to the side condition $\int_{\Omega} q(x, t) d x=0$ for almost all $t$ ).
$2^{\circ}$ When $s \geq \frac{n+2}{p}-3\left[s>\frac{n+2}{p}-3\right.$ if $\left.\frac{n}{2 p}-\frac{3}{2} \in \mathbf{N}_{+}, p \neq 2\right]$, there is a constant $N_{s, p, b}$ such that for data $\Phi_{k}$ with data norm $\mathcal{N}_{s, p, b}^{(k)}\left(\Phi_{k}\right)<N_{s, p, b}$ there exists a solution $\{u, q\}$ of (1.1) with (2.12), the norm depending continuously on $\Phi_{k}$. When $s \geq s_{0}$ for some $s_{0}>\frac{n+2}{p}-3\left[\frac{s_{0}}{2}-\frac{1}{p} \notin \mathbf{N}_{+}\right.$if $\left.p \neq 2\right]$, the norm condition for existence can be replaced by the condition $\mathcal{N}_{s_{0}, p, b}^{(k)}\left(\Phi_{k}\right)<N_{s_{0}, p, b}$.
$3^{\circ}$ When $s>\frac{n+2}{p}-3$, one can for each $N>0$ choose $b^{\prime} \leq b$ such that there exists a solution $\{u, q\}$ of (1.1) satisfying (2.12) with $b$ replaced by $b^{\prime}$, and with norm depending continuously on $\Phi_{k}$, for any set of data $\Phi_{k}$ with norm $\mathcal{N}_{s, p, b^{\prime}}^{(k)}\left(\Phi_{k}\right)<N$. For $s \geq s_{0}, s_{0}$ as above, the solution can be obtained with $b^{\prime}$ defined relative to $s_{0}$.

The statements hold with $H_{p}$ replaced by $B_{p}$ throughout, even without the conditions in [...].

One concludes furthermore that $q \in H_{p}^{(s+1, s / 2)}\left(\bar{Q}_{I_{b^{\prime}}}\right)$ when $s \geq 0$ or $f$ is as in (2.4); in some cases $q$ belongs to a better space, see $\left[{ }^{7}\right]$, Th. 3.6.

For $k=0, s=0$, the result is consistent with Solonnikov's result [ ${ }^{18}$ ], Th. 10.1 for $n=3$, showing the existence of solutions in $W_{p}^{(2,1)}\left(\bar{Q}_{I_{b}}\right)^{n}$ to the Dirichlet Navier-Stokes problem when $f \in L_{p}(\Omega)^{n}, \varphi=0, u_{0} \in$ $\operatorname{pr}_{J_{0}} B_{p}^{2-2 / p}(\bar{\Omega})^{n}$ and $p \geq \frac{5}{3}$.

When both $f$ and $\varphi_{k}$ are 0 , one can get solutions with still more general initial data, e.g. in $L_{n}(\Omega)$, cf. [ ${ }^{18}$ ] for the Dirichlet problem $(n=3)$, GigaMiyakawa [ ${ }^{3}$ ] and Giga [ ${ }^{2}$ ] for Dirichlet and intermediate problems, and [ ${ }^{8}$ ] for Neumann and Dirichlet problems. In [ ${ }^{8}$, we use the semigroup $U(t)$ associated with $A_{k}=\left(-\Delta+G_{k}\right)_{T_{k}^{\prime}}$ to obtain solutions e.g. in spaces $C^{0}\left(I_{b^{\prime}} ; H_{p}^{r}(\bar{\Omega})^{n}\right)$, when $u_{0}$ is taken in $H_{p}^{r}(\bar{\Omega})^{n}$, allowed for $r \geq \frac{n}{p}-1$. [ $\left.{ }^{18}\right]$, [ $\left.{ }^{3}\right]$, [ $\left.{ }^{[ }\right]$and von Wahl $\left[{ }^{20}\right]$ also treat problems with $f \neq 0, \varphi_{0}=0$, in related spaces.

For the Dirichlet problem, we can include infinite intervals $I=\mathbf{R}_{+}$in certain cases:

Theorem 2.2 Hypotheses as in Theorem 2.1. In the Dirichlet case ( $k=0$ ), the existence of solutions with (2.12) for sufficiently small data extends to $b=$
$+\infty$ (generalizing Theorem $2.12^{\circ}$ ), when either $1^{\circ}, 2^{\circ}$ or $3^{\circ}$ holds in addition to the conditions $s>\frac{1}{p}-1, s \geq \frac{n+2}{p}-3$ :
$1^{\circ} n \leq 4$.
$2^{\circ} s<\frac{3}{p}$.
$3^{\circ}$ The data have vanishing initial values, i.e., satisfy

$$
\begin{equation*}
\left\{f, \varphi_{0}, u_{0}\right\} \in H_{p(0)}^{(s, s / 2)}\left(\bar{Q}_{\mathbf{R}_{+}}\right)^{n} \times B_{p(0)}^{\left(s+2-\frac{1}{p},\left(s+2-\frac{1}{p}\right) / 2\right)}\left(\underline{E}_{\tau, \mathbf{R}_{+}}\right) \times\{0\} . \tag{2.13}
\end{equation*}
$$

There is a similar generalization with $H_{p}$ replaced by $B_{p}$.
The method of proof of Theorem 2.1 in $\left[^{7}\right]$ consists of the following four steps: 1) Reduction of the linearized problem to a truly parabolic but pseudodifferential initial-boundary value problem ( $\left.\left.\left[{ }^{11}\right],{ }^{[4}\right]\right)$. 2) Solution of the linear reduced parabolic problem by pseudo-differential machinery (from $\left[{ }^{9}\right],\left[{ }^{6}\right]$ ). 3) Solution of the corresponding reduced nonlinear pseudodifferential problem, by use of product estimates and iteration. 4) Conclusions for the original nonlinear problem. For Theorem 2.2 one uses moreover, that the resolvent of the linearized stationary problem is really only applied to the solenoidal space, where the spectrum for $k=0$ is a closed subset of $\mathbf{R}_{+}$; this allows sharper estimates.

## 3 A simplified method

We shall now explain the method of proof in a version where Step 2 is simplified.
We first treat the associated resolvent problem, where $\partial_{t}$ is replaced by the complex parameter $-\lambda$. To be concrete, consider (1.2) in the Neumann case $k=1$ (which has been studied less than the Dirichlet case $k=0$ ):

$$
\begin{align*}
(-\Delta-\lambda) u+\operatorname{grad} q & =f & & \text { on } \Omega, \\
\operatorname{div} u & =0 & & \text { on } \Omega,  \tag{3.1}\\
\chi_{1} u-\gamma_{0} q & =\varphi & & \text { on } \Gamma,
\end{align*}
$$

with $f$ and $\varphi$ given in $H_{p}^{s}(\bar{\Omega})^{n}$ resp. $B_{p}^{s+1-\frac{1}{p}}(\Gamma)^{n} ; \Omega$ bounded and smooth.
Applying - div to the first line in (3.1) and taking the normal component of the third line, we find:

$$
\begin{align*}
-\Delta q & =-\operatorname{div} f \\
\gamma_{0} q & =2 \gamma_{1} u_{\nu}-\varphi_{\nu} \tag{3.2}
\end{align*}
$$

This is a Dirichlet problem for $q$, so if we denote $\binom{-\Delta}{\gamma_{0}}^{-1}=\left(R_{D} K_{D}\right)$, we have

$$
\begin{gather*}
q=-R_{D} \operatorname{div} f+K_{D}\left(2 \gamma_{1} u_{\nu}-\varphi_{\nu}\right), \text { where } \\
R_{D}: H_{p}^{s-1}(\bar{\Omega}) \rightarrow H_{p}^{s+1}(\bar{\Omega}) \text { for } s>\frac{1}{p}-1,  \tag{3.3}\\
K_{D}: B_{p}^{s+1-\frac{1}{p}}(\Gamma) \rightarrow H_{p}^{s+1}(\bar{\Omega}) \text { for } s \in \mathbf{R} .
\end{gather*}
$$

Insertion of $q$ into (3.1) gives the equations for $u$ :

$$
\begin{align*}
(-\Delta-\lambda) u+2 \operatorname{grad} K_{D} \gamma_{1} \operatorname{pr}_{\nu} u & =f+\operatorname{grad} R_{D} \operatorname{div} f+\operatorname{grad} K_{D} \varphi_{\nu} \\
\gamma_{0} \operatorname{div} u & =0 \\
\operatorname{pr}_{\tau} \chi_{1} u & =\varphi_{\tau} \tag{3.4}
\end{align*}
$$

A solution of $u$ of (3.4) will satisfy (3.1) when $q$ is defined from $u, f$ and $\varphi$ by (3.3) (note that div grad $K_{D}=0$ ).

We write (3.4) in the short form

$$
\begin{align*}
(-\Delta-\lambda) u+K T u & =f_{1} & & \text { on } \Omega \\
T_{1}^{\prime} u & =\psi_{1} & & \text { on } \Gamma \tag{3.5}
\end{align*}
$$

where we have set

$$
\begin{gather*}
K=2 \operatorname{grad} K_{D}, \quad T=\gamma_{1} \operatorname{pr}_{\nu}, \quad T_{1}^{\prime}=\left\{\operatorname{pr}_{\tau} \chi_{1}, \gamma_{0} \operatorname{div}\right\} \\
f_{1}=f+\operatorname{grad} R_{D} \operatorname{div} f+\operatorname{grad} K_{D} \varphi_{\nu}, \quad \psi_{1}=\left\{\varphi_{\tau}, 0\right\} \tag{3.6}
\end{gather*}
$$

here $I+\operatorname{grad} R_{D}$ div equals the projection operator $\operatorname{pr}_{J}$ that maps $H_{p}^{s}(\bar{\Omega})^{n}$ onto the solenoidal space $J_{p}^{s}=\left\{u \in H_{p}^{s}(\bar{\Omega})^{n} \mid \operatorname{div} u=0\right\}$ for $s>\frac{1}{p}-1$ (cf. [ $\left.{ }^{5}\right]$, Example 3.14); $K$ is a Poisson operator of order 1 , and $T$ and $T_{1}^{\prime}$ are trace operators of order 1.

In order to use other known properties of the Laplace operator, we now make a new reduction. Write the problem (3.5) as follows:

$$
\begin{align*}
(-\Delta-\lambda) u & =f_{1}-K T u & & \text { on } \Omega \\
\gamma_{1} u & =\psi_{1}-\left(T_{1}^{\prime}-\gamma_{1}\right) u & & \text { on } \Gamma . \tag{3.7}
\end{align*}
$$

The system $\left\{-\Delta-\lambda, \gamma_{1}\right\}$ is uniformly parameter-elliptic (in the sense of $\left[{ }^{9}\right]$, $\left.\left.{ }^{[6}\right]\right)$ for $\lambda$ on rays with argument $\left.\theta \in\right] 0,2 \pi\left[\right.$; and it is bijective for $\lambda \in \mathbf{C} \backslash \overline{\mathbf{R}}_{+}$. By a simple application of $\left[{ }^{6}\right]$, the inverse is continuous for each $s>\frac{1}{p}-1$,

$$
\begin{equation*}
\binom{-\Delta-\lambda}{\gamma_{1}}^{-1}=\left(R_{N, \lambda} K_{N, \lambda}\right): H_{p}^{s, \mu}(\bar{\Omega})^{n} \times B_{p}^{s+1-\frac{1}{p}, \mu}(\Gamma)^{n} \rightarrow H_{p}^{s+2, \mu}(\bar{\Omega})^{n} \tag{3.8}
\end{equation*}
$$

uniformly for $\lambda$ in sets $V_{\varepsilon}, \varepsilon>0$,

$$
\begin{equation*}
V_{\varepsilon}=\{\lambda \in \mathbf{C}|\arg \lambda \in[\varepsilon, 2 \pi-\varepsilon],|\lambda| \geq \varepsilon\} \tag{3.9}
\end{equation*}
$$

here $\mu=|\lambda|^{\frac{1}{2}}$. (The $H_{p}^{s, \mu}$ and $B_{p}^{s, \mu}$ spaces are $H_{p}^{s}$ and $B_{p}^{s}$ spaces provided with norms depending on $\mu$, as in the basic case of $H_{p}^{s, \mu}\left(\mathbf{R}^{n}\right)$, which is the space provided with the norm $\left\|\mathrm{OP}\left(\left(|\xi|^{2}+|\mu|^{2}+1\right)^{s / 2}\right) u\right\|_{p}$, cf. $\left[{ }^{9}\right]$. Mapping properties like (3.8) are well-known in the literature, cf. e.g. [ ${ }^{15}$ ], except perhaps for the extension to low values of $s$.) With (3.8), we can write (3.7) as:

$$
\begin{equation*}
u=R_{N, \lambda}\left(f_{1}-K T u\right)+K_{N, \lambda}\left(\psi_{1}-\left(T_{1}^{\prime}-\gamma_{1}\right) u\right) \tag{3.10}
\end{equation*}
$$

or, if we set

$$
\begin{equation*}
\Phi=R_{N, \lambda} f_{1}+K_{N, \lambda} \psi_{1}, \quad \mathcal{K}_{\lambda}=\left(R_{N, \lambda} K K_{N, \lambda}\right), \quad \mathcal{T}=\binom{T}{T_{1}^{\prime}-\gamma_{1}} \tag{3.11}
\end{equation*}
$$

as:

$$
\begin{equation*}
\left(I+\mathcal{K}_{\lambda} \mathcal{T}\right) u=\Phi \tag{3.12}
\end{equation*}
$$

We observe that the operators in (3.11), when considered as depending on the parameter $\lambda$, have regularity $\frac{1}{2}$ in the sense of [ ${ }^{4}$, since $R_{N, \lambda}, K_{N, \lambda}$ and $\mathcal{T}$ have regularity $+\infty$, and $K$, being of order 1 , counts with regularity $\frac{1}{2}$ by [ ${ }^{4}$ ], Prop. 2.3.14.

Now we need the elementary

Lemma 3.1 Let $A: V \rightarrow W$ and $B: W \rightarrow V$ be linear mappings. If $I+A B:$ $W \rightarrow W$ is bijective, then $I+B A: V \rightarrow V$ is bijective, with

$$
\begin{equation*}
(I+B A)^{-1}=I-B(I+A B)^{-1} A \tag{3.13}
\end{equation*}
$$

Proof: One just has to check:

$$
\begin{align*}
& (I+B A)\left(I-B(I+A B)^{-1} A\right) \\
& \quad=I+B A-B(I+A B)^{-1} A-B A B(I+A B)^{-1} A \\
& \quad=I+B A-B(I+A B)(I+A B)^{-1} A=I \tag{3.14}
\end{align*}
$$

with a similar calculation for the left composition.

The lemma will be applied with

$$
\begin{align*}
& A=\mathcal{K}_{\lambda}: B_{p}^{s+1-\frac{1}{p}, \mu}(\Gamma)^{n+1} \rightarrow H_{p}^{s+2, \mu}(\bar{\Omega})^{n} \\
& B=\mathcal{T}: H_{p}^{s+2, \mu}(\bar{\Omega})^{n} \rightarrow B_{p}^{s+1-\frac{1}{p}, \mu}(\Gamma)^{n+1}, \quad s>\frac{1}{p}-1, \quad \lambda \in V_{\varepsilon} \tag{3.15}
\end{align*}
$$

and also with the roles of $A$ and $B$ interchanged. The lemma shows that (21) can be uniquely solved (in these spaces) if and only if $I+\mathcal{T} \mathcal{K}_{\lambda}$ is invertible, in which case

$$
\begin{equation*}
\left(I+\mathcal{K}_{\lambda} \mathcal{T}\right)^{-1}=I-\mathcal{K}_{\lambda}\left(I+\mathcal{T} \mathcal{K}_{\lambda}\right)^{-1} \mathcal{T} \tag{3.16}
\end{equation*}
$$

Now $I+\mathcal{T} \mathcal{K}_{\lambda}$ is much easier to deal with than $I+\mathcal{K}_{\lambda} \mathcal{T}$, since it is a parameter-dependent ps.d.o. on the boundaryless compact manifold $\Gamma$ ! In details, $\mathcal{T} \mathcal{K}_{\lambda}$ is an $(n+1) \times(n+1)$-matrix

$$
\mathcal{T} \mathcal{K}_{\lambda}=\left(\begin{array}{cc}
T R_{N, \lambda} K & T K_{N, \lambda}  \tag{3.17}\\
\left(T_{1}^{\prime}-\gamma_{1}\right) R_{N, \lambda} K & \left(T_{1}^{\prime}-\gamma_{1}\right) K_{N, \lambda}
\end{array}\right)
$$

where the entries are of regularity $\frac{1}{2}$, in the sense of [ ${ }^{4}$ ].
It is parameter-elliptic in the sense of [ ${ }^{4}$ ], for $\lambda$ on rays in $\mathbf{C} \backslash \mathbf{R}_{+}$, since this is a question of bijectiveness of certain model operators at the boundary (and certain matrices), a property that can be traced all the way from (3.4) to (3.12); the parameter-ellipticity of (3.4) was shown in $\left[{ }^{14}\right]$, Sect. 6. This implies that for any $\varepsilon>0$, there is an $r(\varepsilon)>0$ so that for any $t \in \mathbf{R}$,

$$
\begin{gather*}
I+\mathcal{T} \mathcal{K}_{\lambda}: B_{p}^{t, \mu}(\Gamma)^{n+1} \underset{\rightarrow}{\rightarrow} B_{p}^{t, \mu}(\Gamma)^{n+1}, \text { uniformly for } \lambda \in W_{\varepsilon, r(\varepsilon)} \\
W_{\varepsilon, r(\varepsilon)}=\{\lambda \in \mathbf{C}|\arg \lambda \in[\varepsilon, 2 \pi-\varepsilon],|\lambda| \geq r(\varepsilon)\} \tag{3.18}
\end{gather*}
$$

we denote the inverse

$$
\begin{equation*}
\left(I+\mathcal{T} \mathcal{K}_{\lambda}\right)^{-1}=Q_{\lambda} \tag{3.19}
\end{equation*}
$$

For such $\lambda$ we also have the inverse, by Lemma 3.1,

$$
\begin{equation*}
\left(I+\mathcal{K}_{\lambda} \mathcal{T}\right)^{-1}=I-\mathcal{K}_{\lambda} Q_{\lambda} \mathcal{T}: H_{p}^{s+2, \mu}(\Omega)^{n} \stackrel{\sim}{\rightarrow} H_{p}^{s+2, \mu}(\Omega)^{n}, \tag{3.20}
\end{equation*}
$$

uniformly for $\lambda \in W_{\varepsilon, r(\varepsilon)}$, when $s>\frac{1}{p}-1$.
Altogether, we solve (3.1) by taking

$$
\begin{align*}
u= & \left(I+\mathcal{K}_{\lambda} \mathcal{T}\right)^{-1} \Phi=\left(I-\mathcal{K}_{\lambda} Q_{\lambda} \mathcal{T}\right)\left(R_{N, \lambda} f_{1}+K_{N, \lambda} \varphi_{1}\right) \\
= & \left(I-\mathcal{K}_{\lambda} Q_{\lambda} \mathcal{T}\right)\left(R_{N, \lambda}\left(f+\operatorname{grad} R_{D} \operatorname{div} f+\operatorname{grad} K_{D} \varphi_{\nu}\right)\right. \\
& \left.\quad+K_{N, \lambda}\left\{\varphi_{\tau}, 0\right\}\right) \tag{3.21}
\end{align*}
$$

and defining $q$ by (3.3). The point is here that all operators except the factor $Q_{\lambda}$ stem from classical resolvent problems for the Laplace operator; and $Q_{\lambda}$ is a parameter-dependent ps.d.o. on $\Gamma$.

For the other boundary conditions (the cases $k=0,2,3,4$ in (1.2)), there are similar methods; for $k=0,2,4$, the roles of the Dirichlet and Neumann problems are interchanged.

Also in the original problem where $-\lambda$ is replaced by $\partial_{t}$, this approach gives some simplifications. Indeed, as described in $\left[{ }^{7}\right]$, the resolvent considerations carry over to solvability of the $t$-dependent problem with initial data 0 , formally by a Laplace transformation. Analogously to the derivation of (3.21) from (3.5) we find that the problem

$$
\begin{align*}
\left(\partial_{t}-\Delta-\varrho\right) v+K T v & =g_{1} & & \text { on } Q_{\mathbf{R}}  \tag{3.22}\\
T_{1}^{\prime} v & =\zeta_{1} & & \text { on } S_{\mathbf{R}}
\end{align*}
$$

with $v(x, t), g_{1}(x, t)$ and $\zeta_{1}(x, t)$ supported for $t \geq 0$, and $\varrho<\inf \{\operatorname{Re} \lambda \mid \lambda \in$ $\left.W_{\varepsilon, r(\varepsilon)}\right\}$, has a solution operator described by

$$
\begin{equation*}
v=(I-\mathbf{K Q} \mathcal{T})\left(\mathbf{R}_{\mathbf{N}} g_{1}+\mathbf{K}_{\mathbf{N}} \zeta_{1}\right) \tag{3.23}
\end{equation*}
$$

with

$$
\begin{align*}
\left(\mathbf{R}_{\mathbf{N}} \mathbf{K}_{\mathbf{N}}\right) & =\binom{\partial_{t}-\Delta-\varrho}{\gamma_{1}}^{-1} \\
\mathbf{K} & =\left(\mathbf{R}_{\mathbf{N}} K \mathbf{K}_{\mathbf{N}}\right)  \tag{3.24}\\
\mathbf{Q} & =(I+\mathcal{T} \mathbf{K})^{-1}
\end{align*}
$$

The latter exists since $I+\mathcal{T} \mathbf{K}$ is derived from the parameter-dependent operator $I+\mathcal{T} \mathcal{K}_{\lambda+\varrho}$ by replacing $-\lambda$ by $i \tau$ in the symbol and using a pseudodifferential definition in one more variable; here $I+\mathcal{T} \mathcal{K}_{\lambda+\varrho}$ is invertible for $\lambda$ in an obtuse neighborhood of $\{\operatorname{Re} \lambda \leq 0\}$, and the calculus of $\left[{ }^{6}\right]$, Theorem 3.1 $1^{\circ}$ is applicable.

The resulting estimates for the solution operator are the same as those described in detail in $\left[^{7}\right]$. It is because of the constant $\varrho$ in (3.22), that we do not obtain time-global estimates in $H_{p}^{(s, s / 2)}$ spaces in general. See however the special considerations for the Dirichlet problem in [ ${ }^{7}$ ].

## 4 Exterior problems

Consider now the case where $\Omega$ is the complement of a compact set in $\mathbf{R}^{n}$, still with smooth boundary $\Gamma$. We can then investigate how the method of Section

3 can be used. Applying - div to the first line in (3.1) and taking the normal component of the third line, we again arrive at (3.2), now an exterior Dirichlet problem for $q$.

Theorem 4.1 The exterior Dirichlet problem

$$
\begin{equation*}
-\Delta v=g \text { in } \Omega, \quad \gamma_{0} v=\psi \text { on } \Gamma \tag{4.1}
\end{equation*}
$$

has a solution operator $\left(R_{D} K_{D}\right)$ such that

$$
\begin{gather*}
\operatorname{grad} K_{D}: B_{p}^{s+1-\frac{1}{p}}(\Gamma) \rightarrow H_{p}^{s}(\bar{\Omega})^{n} \text { for } s \in \mathbf{R} \\
\operatorname{grad} R_{D} \operatorname{div}: H_{p}^{s}(\bar{\Omega})^{n} \rightarrow H_{p}^{s}(\bar{\Omega})^{n} \text { for } s>\frac{1}{p}-1 \tag{4.2}
\end{gather*}
$$

and $\operatorname{grad} K_{D}$ maps into $\bigcap_{q>1, r \in \mathbf{R}} H_{q}^{r}(\{|x| \geq R\}$ ), when $\{|x| \geq R\} \subset \Omega$. (More precisely, $R_{D}$ is defined for $g$ with compact support, and $\operatorname{grad} R_{D}$ div is extended by continuity.)

Here $K_{D}$ is uniquely determined by the property that grad $K_{D}$ should map into functions that are $O\left(|x|^{-n}\right)$ for $|x| \rightarrow \infty$. It is also uniquely determined by requiring grad $K_{D}$ to map into $\bigcap_{q>1} L_{q}(\{|x| \geq R\})$ - or just into $L_{p}(\{|x| \geq$ $R\}$ ), if $p \leq n /(n-1)$.

Proof: The mapping $K_{D}$ is constructed as follows: We want to find a solution of

$$
\begin{equation*}
-\Delta v=0 \text { in } \Omega, \quad \gamma_{0} v=\psi \text { on } \Gamma \tag{4.3}
\end{equation*}
$$

for given $\psi \in B_{p}^{s+1-\frac{1}{p}}(\Gamma)$, such that $\operatorname{grad} v$ is $O\left(|x|^{-n}\right)$ for $|x| \rightarrow \infty$. (The derivatives will also be $O\left(|x|^{-n}\right)$, and then $\operatorname{grad} u \in H_{q}^{r}(\{|x| \geq R\})$ for all $r$, all $q>1$, when $\{|x| \geq R\} \subset \Omega$.) Instead of (4.3), we can study

$$
\begin{equation*}
-\Delta v_{c}=0 \text { in } \Omega, \quad \gamma_{0} v_{c}=\psi-c \text { on } \Gamma \tag{4.4}
\end{equation*}
$$

where $c$ is a constant to be chosen freely; if $v_{c}$ solves (4.4), then $v=v_{c}+c$ solves (4.3) (and vice versa), and they have the same gradient.

We can assume that 0 is in the complement of $\bar{\Omega}$, so that the inversion $x \mapsto x /|x|^{2}$ maps $\Omega$ onto $\Omega^{*} \backslash\{0\}$, where $\Omega^{*}$ is a bounded open smooth set with $0 \in \Omega^{*}$, and $\partial \Omega^{*}=\Gamma^{*}$ is the image of $\Gamma$. (In the following, let $n \geq 3$; there is a similar proof for $n=2$.) Let $W(x)$ be the solution of the special Dirichlet problem

$$
\begin{equation*}
-\Delta W=0 \text { in } \Omega^{*}, \quad \gamma_{0} W(x)=|x|^{2-n} \text { on } \Gamma^{*} ; \tag{4.5}
\end{equation*}
$$

since $|x|^{2-n}>0$ on $\Gamma^{*}, W(0)>0$ by the maximum principle.
Now if a function $V$ is harmonic in $\Omega^{*} \backslash\{0\}$, then the function $|x|^{2-n} V\left(x /|x|^{2}\right)$ is harmonic in $\Omega$, and vice versa (the Kelvin transformation). Since $\psi$ on $\Gamma$ is carried over to $|x|^{2-n} \psi\left(x /|x|^{2}\right)$ on $\Gamma^{*}$, lying in $B_{p}^{s+1-\frac{1}{p}}\left(\Gamma^{*}\right)$, we can find a solution of (4.4) by solving

$$
\begin{equation*}
-\Delta V_{c}=0 \text { in } \Omega^{*}, \quad \gamma_{0} V_{c}=|x|^{2-n}\left(\psi\left(x /|x|^{2}\right)-c\right) \text { on } \Gamma^{*} \tag{4.6}
\end{equation*}
$$

The problem (4.6) has a unique solution $V_{c} \in H_{p}^{s+1}\left(\bar{\Omega}^{*}\right) \cap C^{\infty}\left(\Omega^{*}\right)$ for each $c$; and by the linearity, $V_{c}=V_{0}-c W$, cf. (4.5). Take

$$
\begin{equation*}
c^{*}=V_{0}(0) / W(0) \tag{4.7}
\end{equation*}
$$

then $V_{c^{*}}(0)=0$. Now let

$$
\begin{equation*}
v_{c^{*}}(x)=|x|^{2-n} V_{c^{*}}\left(x /|x|^{2}\right) \tag{4.8}
\end{equation*}
$$

it solves (4.4). Since $V_{c^{*}}(0)=0, D^{\alpha} v_{c^{*}}(x)$ is $O\left(|x|^{1-n-|\alpha|}\right)$ for $|x| \rightarrow \infty$, any $\alpha$ (seen from the Taylor expansion of $V_{c^{*}}$ at 0 ). In particular, grad $v_{c^{*}}$ and its derivatives are $O\left(|x|^{-n}\right)$ and hence $L_{q}$ integrable at $\infty$ for any $q>1$. On the other hand, $V_{c^{*}} \in H_{p}^{s+1}\left(\bar{\Omega}^{*}\right)$ implies that $v_{c^{*}}$ is in $H_{p}^{s+1}$ over bounded subsets of $\bar{\Omega}$. Altogether, it is found that $v_{c^{*}}$, and hence $v=v_{c^{*}}+c^{*}$, have gradient in $H_{p}^{s}(\bar{\Omega})$ and in $\bigcap_{q>1, r \in \mathbf{R}} H_{q}^{r}(\{|x| \geq R\})$. Defining $K_{D}$ as the mapping from $\psi$ to $v$, we have obtained an operator with the asserted mapping properties.

To show the uniqueness, let $u$ be a solution of (4.3) with $\operatorname{grad} u=O\left(|x|^{-n}\right)$ for $|x| \rightarrow \infty$. Recall that any function $u(x)$ that is harmonic on $\Omega$ has a unique Laurent expansion

$$
\begin{equation*}
u(x)=\sum_{k=0}^{\infty} H_{k}^{*}(x)+\sum_{k=0}^{\infty} \frac{H_{k}(x)}{|x|^{n-2+2 k}} \tag{4.9}
\end{equation*}
$$

where the functions $H_{k}^{*}(x)$ and $H_{k}(x)$ are homogeneous harmonic polynomials of degree $k$; cf. Brelot [ ${ }^{1}$, pp. 197-202, where also sets $\left\{R_{1}<|x|<R_{2}\right\}$ are considered. The first series in (4.9) converges uniformly on bounded sets, the second converges uniformly on sets $\{|x| \geq R\}$ contained in $\Omega$, and the derivatives of $u$ are represented by the termwise differentiated expressions. When $u$ has the form (4.9),

$$
\begin{equation*}
\operatorname{grad} u(x)=\sum_{k=1}^{\infty} \operatorname{grad} H_{k}^{*}(x)+\sum_{k=0}^{\infty} \operatorname{grad} \frac{H_{k}(x)}{|x|^{n-2+2 k}} \tag{4.10}
\end{equation*}
$$

where $\operatorname{grad}\left(H_{k}(x) /|x|^{n-2+2 k}\right)$ is $O\left(|x|^{1-n-k}\right)$ for each $k \geq 0$. The requirement that $\operatorname{grad} u$ should be $O\left(|x|^{-n}\right)$ rules out the polynomials $H_{k}^{*}$ for $k \geq 1$ as well as the term $H_{0} /|x|^{n-2}$, so

$$
\begin{equation*}
u(x)=H_{0}^{*}+v(x), \quad v(x)=\sum_{k=1}^{\infty} \frac{H_{k}(x)}{|x|^{n-2+2 k}} \tag{4.11}
\end{equation*}
$$

Now the Kelvin transform of $v, V(x)=|x|^{2-n} v\left(x /|x|^{2}\right)$ has the representation

$$
\begin{equation*}
V(x)=\sum_{k=1}^{\infty} H_{k}(x) \tag{4.12}
\end{equation*}
$$

it is $C^{\infty}$ on $\Omega^{*}$ with $V(0)=0$. Thus $V$ is the unique solution of (33) with $c=H_{0}^{*}$. If $\widetilde{u}=\widetilde{c}+\widetilde{v}(x)$ also solves (4.3), with $\widetilde{v}=\sum_{k=1}^{\infty} \widetilde{H}_{k}(x) /|x|^{n-2+2 k}$, then $\widetilde{v}$ is the Kelvin transform of the unique solution $\widetilde{V}$ of (4.6) with $c=\widetilde{c}$; this also has $\tilde{V}(0)=0$. Now $V-\widetilde{V}=\left(-H_{0}^{*}+\widetilde{c}\right) W$ (cf. (4.5)), and since $V(0)=\widetilde{V}(0)=0$ and $W(0) \neq 0, \widetilde{c}$ must equal $H_{0}^{*}$, and $V=\widetilde{V}$. Thus $u$ is uniquely determined.

To require $\operatorname{grad} u \in L_{q}(\{|x| \geq R\})$ for some $q \leq n /(n-1)$ likewise reduces $u$ to the form (4.11), since nonzero polynomials and $|x|^{1-n} \notin L_{q}(\{|x| \geq R\})$; and the analysis goes as above.

Now consider $R_{D}$. To solve the problem (4.1) with $\psi=0$ and a nonzero $g=\operatorname{div} h, h \in H_{p}^{s}(\bar{\Omega})^{n}$, we let $l$ be a continuous linear extension operator $l: H_{p}^{s}(\bar{\Omega}) \rightarrow H^{s}\left(\mathbf{R}^{n}\right)$ such that $\operatorname{div} l f=l \operatorname{div} f$, and search for $v$ in the form $v=r_{\Omega} v_{1}+v_{2}$, where $v_{1}$ and $v_{2}$ solve

$$
-\Delta v_{1}=l g \text { on } \mathbf{R}^{n}, \text { resp. }\left\{\begin{array}{l}
-\Delta v_{2}=0 \text { in } \Omega  \tag{4.13}\\
\gamma_{0} v_{2}=-\gamma_{0} v_{1} \text { on } \Gamma
\end{array}\right.
$$

( $r_{\Omega}$ denotes restriction to $\Omega$.) When $g=\operatorname{div} h$, we want $\operatorname{grad} v$ to depend continuously on $h$ in $H_{p}^{s}(\bar{\Omega})^{n}$.

To begin with, let $g$ have compact support. To get a convenient solution of the problem for $v_{1}$, we define $R$ as the operator (for compactly supported f)

$$
\begin{gather*}
R: f \mapsto \omega_{n}|x|^{2-n} * f-c(f) \\
\text { where } c(f) \text { is the constant } \frac{\int_{\Gamma} \gamma_{0}\left(\omega_{n}|x|^{2-n} * f\right) d \sigma}{\int_{\Gamma} 1 d \sigma} \tag{4.14}
\end{gather*}
$$

and $\omega_{n}|x|^{2-n}$ is the Newton potential $\left(\omega_{n}|x|^{2-n} * f\right.$ can be viewed as $\left.\mathrm{OP}\left(|\xi|^{-2}\right) f\right)$. This satisfies $-\Delta R f=f$, and

$$
\begin{equation*}
\int_{\Gamma} \gamma_{0} R f d \sigma=0 . \tag{4.15}
\end{equation*}
$$

When $g \in H_{p}^{s-1}(\bar{\Omega})$, then $v_{1}=R l g \in H_{p, \text { loc }}^{s+1}\left(\mathbf{R}^{n}\right)$ by elliptic regularity, so $\gamma_{0} v_{1} \in B_{p}^{s+1-\frac{1}{p}}(\Gamma) \subset L_{p}(\Gamma)$, and the expressions are well-defined.

Now insert $\gamma_{0} v_{1}$ in the equations for $v_{2}$; then the operator $K_{D}$ established above gives a solution $v_{2}$ with grad $v_{2} \in H_{p}^{s}(\bar{\Omega})^{n}$. Altogether, we take

$$
\begin{equation*}
R_{D} g \equiv v=r_{\Omega} v_{1}+v_{2}=r_{\Omega} R l g-K_{D} \gamma_{0} R l g . \tag{4.16}
\end{equation*}
$$

In particular, when $g=\operatorname{div} h$ is inserted, then

$$
\begin{equation*}
\operatorname{grad} v=r_{\Omega} \operatorname{grad} v_{1}+\operatorname{grad} v_{2}=r_{\Omega} \operatorname{grad} R \operatorname{div} l h-\operatorname{grad} K_{D} \gamma_{0} R \operatorname{div} l h . \tag{4.17}
\end{equation*}
$$

Here we have for grad $v_{1}$, by Fourier transformation,

$$
\begin{equation*}
\operatorname{grad} v_{1}=\operatorname{grad} R \operatorname{div} l h=\mathrm{OP}\left(\left(\xi_{i} \xi_{j} /|\xi|^{2}\right)_{i, j=1, \ldots, n}\right) l h \tag{4.18}
\end{equation*}
$$

The operator $\operatorname{OP}\left(\left(\xi_{i} \xi_{j} /|\xi|^{2}\right)_{i, j=1, \ldots, n}\right)$ extends to a continuous operator in $L_{p}\left(\mathbf{R}^{n}\right)^{n}$ by a result of Calderón and Zygmund; and it is likewise continuous in $H_{p}^{s}\left(\mathbf{R}^{n}\right)^{n}$ for all $s \in \mathbf{R}$, since

$$
\begin{equation*}
\mathrm{OP}\left(\langle\xi\rangle^{s}\right) \mathrm{OP}\left(\left(\xi_{i} \xi_{j} /|\xi|^{2}\right)_{i, j=1, \ldots, n}\right) \mathrm{OP}\left(\langle\xi\rangle^{-s}\right)=\mathrm{OP}\left(\left(\xi_{i} \xi_{j} /|\xi|^{2}\right)_{i, j=1, \ldots, n}\right) . \tag{4.19}
\end{equation*}
$$

Thus

$$
\begin{equation*}
\left\|r_{\Omega} \operatorname{grad} v_{1}\right\|_{H_{p}^{s}(\bar{\Omega})} \leq\left\|\operatorname{grad} v_{1}\right\|_{H_{p}^{s}\left(\mathbf{R}^{n}\right)} \leq C_{1}\|l h\|_{H_{p}^{s}\left(\mathbf{R}^{n}\right)} \leq C_{2}\|h\|_{H_{p}^{s}(\bar{\Omega})} . \tag{4.20}
\end{equation*}
$$

This extends to $h$ with arbitrary support, by approximation by $\eta(x / N) h$, for $\eta \in C_{0}^{\infty}\left(\mathbf{R}^{n}\right)$ equal to 1 near 0 and $N \rightarrow \infty$.

For $v_{2}$, we need to know that the mapping from $h$ (or just from grad $v_{1}$ ) to $\gamma_{0} v_{1}$ is continuous from $H_{p}^{s}(\bar{\Omega})$ to $B_{p}^{s+1-\frac{1}{p}}(\Gamma)$. This is shown by the help of Lemma 4.2 below. Take for $\Xi$ a large ball containing $\Gamma$ in its interior. In view of (4.15), the lemma gives that

$$
\begin{align*}
\left\|\gamma_{0} v_{1}\right\|_{B_{p}^{s+1-p}(\Gamma)} & \leq C_{3}\left\|v_{1}\right\|_{H_{p}^{s+1}(\bar{\Xi})} \leq C C_{3}\left\|\operatorname{grad} v_{1}\right\|_{H_{p}^{s}(\bar{\Xi})} \\
& \leq C C_{3} C_{2}\|h\|_{H_{p}^{s}(\bar{\Omega})} . \tag{4.21}
\end{align*}
$$

Altogether, the mapping $\operatorname{grad} R_{D}$ div, as realized in (4.17), extends by closure to a mapping with the continuity in (4.2).

In the course of the proof we used the following variant of the Poincaré inequality:

Lemma 4.2 Let $s>\frac{1}{p}-1$, and let $\Xi$ be a bounded connected open subset of $\mathbf{R}^{n}$ with a sufficiently smooth boundary such that the injection of $H_{p}^{s+1}(\overline{\bar{\Xi}})$ into $H_{p}^{s}(\bar{\Xi})$ is compact. Let $\Gamma$ be a nonempty smooth closed hypersurface in $\bar{\Xi}$. There is a constant $C$ such that for $u \in H_{p}^{s+1}(\bar{\Xi})$,

$$
\begin{equation*}
\|u\|_{H_{p}^{s}(\bar{\Xi})} \leq C\left(\|\operatorname{grad} u\|_{H_{p}^{s}(\bar{\Xi})}+\left|\int_{\Gamma} \gamma_{0} u d \sigma\right|\right) \tag{4.22}
\end{equation*}
$$

Proof: Note that $\gamma_{0} u=\left.u\right|_{\Gamma}$ is well-defined as an element of $B_{p}^{s+1-\frac{1}{p}}(\Gamma) \subset$ $L_{p}(\Gamma)$, so that the integral has a sense. If a constant $C$ cannot be found, there is a sequence $u_{k}$ with $\left\|u_{k}\right\|_{H_{p}^{s}(\bar{\Xi})}=1$ but $\left\|\operatorname{grad} u_{k}\right\|_{H_{p}^{s}(\bar{\Xi})} \rightarrow 0$ and $\int_{\Gamma} \gamma_{0} u_{k} d \sigma \rightarrow 0$. Since the sequence is bounded in $H_{p}^{s+1}(\bar{\Xi})$, it has a subsequence $u_{k_{j}}$ that is convergent in $H_{p}^{s}(\bar{\Xi})$ to a limit $u_{0}$. Since $\left\|\operatorname{grad} u_{k}\right\|_{H_{p}^{s}(\bar{\Xi})} \rightarrow 0$, the subsequence $u_{k_{j}}$ is convergent in $H_{p}^{s+1}(\bar{\Xi})$, with limit $u_{0}$, and $\gamma_{0} u_{k_{j}} \rightarrow \gamma_{0} u_{0}$ in $B_{p}^{s+1-\frac{1}{p}}(\Gamma)$. Now on one hand $\left\|u_{0}\right\|_{H_{p}^{s}(\bar{\Xi})}=1$, so $u_{0} \neq 0$; on the other hand, $\left\|\operatorname{grad} u_{0}\right\|_{H_{p}^{s}(\bar{\Xi})}=0$ so $u_{0}$ is a constant, and this constant must equal 0 since $\int_{\Gamma} \gamma_{0} u d \sigma=0$. This contradiction proves the statement.

One can also replace the integral in the right hand side of (4.22) by another supplementary term, e.g. the integral of $u$ over some small subdomain.

The mapping $K_{D}$ is established in Simader and Sohr $\left[{ }^{16}\right]$ for integer $s \geq 0$ in a slightly different way, and with lower smoothness assumptions on $\Gamma$; our presentation here was inspired by conversations with B. Fuglede.

The operator $I+\operatorname{grad} R_{D}$ div is the projection onto the solenoidal space $J_{p}^{s}=\left\{u \in H_{p}^{s}(\bar{\Omega})^{n} \mid \operatorname{div} u=0\right\}, s>\frac{1}{p}-1$, generalizing the situation where $\Omega$ is bounded.

Insertion of the formula for $q$ into (3.1) leads to (3.4) and hence (3.5) and (3.7), with the same notation as before. We find by use of Theorem 4.1 that $f_{1} \in H_{p}^{s}(\bar{\Omega})^{n}$, when $f \in H_{p}^{s}(\bar{\Omega})^{n}$ and $\varphi \in B_{p}^{s+1-\frac{1}{p}}(\Gamma)^{n}$.

For exterior domains we also have (3.8), by $\left[{ }^{9}\right]$, so we can write the problem in the form (3.12). Again we can apply Lemma 3.1, using that (3.15) is valid.

It is at this point that it is a great advantage that we can reduce to the inversion of $I+\mathcal{T} \mathcal{K}_{\lambda}$. For this is a ps.d.o. on the compact manifold $\Gamma$. It satisfies the symbol requirements for being parameter-elliptic of regularity $\frac{1}{2}$ on the rays in $\mathbf{C} \backslash \mathbf{R}_{+}$, hence is invertible on these rays for sufficiently large $|\lambda|$, i.e. (3.18) holds. The inverse satisfies (3.20), and we get the solution as in (3.21),

$$
\begin{equation*}
u=\left(I-\mathcal{K}_{\lambda} Q_{\lambda} \mathcal{T}\right)\left(R_{N, \lambda}\left(f+\operatorname{grad} R_{D} \operatorname{div} f+\operatorname{grad} K_{D} \varphi_{\nu}\right)+K_{N, \lambda}\left\{\varphi_{\tau}, 0\right\}\right) \tag{4.23}
\end{equation*}
$$

which lies in $H^{s+2, \mu}(\bar{\Omega})$. Also the considerations for the problem with $-\lambda$ replaced by $\partial_{t}$ go through, and the discussion for nonzero initial values can be completed as in $\left[^{7}\right]$. For the pressure $q$ we use the formula $\operatorname{grad} q=$ $-\operatorname{grad} R_{D} \operatorname{div} f+\operatorname{grad} K_{D}\left(2 \gamma_{1} u_{\nu}-\varphi_{\nu}\right)$, plus the fact that (4.2) carries over to anisotropic spaces as continuous mappings

$$
\begin{align*}
\operatorname{grad} K_{D} & : B_{p}^{\left(s+1-\frac{1}{p},\left(s+1-\frac{1}{p}\right) / 2\right)}\left(\bar{S}_{\mathbf{R}}\right) \rightarrow H_{p}^{(s,(s) / 2)}\left(\bar{Q}_{\mathbf{R}}\right)^{n} \\
\operatorname{grad} R_{D} \operatorname{div} & : H_{p}^{(s, s / 2)}\left(\bar{Q}_{\mathbf{R}}\right)^{n} \rightarrow H_{p}^{(s, s / 2)}\left(\bar{Q}_{\mathbf{R}}\right)^{n} \tag{4.24}
\end{align*}
$$

for $s>\frac{1}{p}-1$. (For $\operatorname{grad} K_{D}$, one uses that the mapping property holds for bounded neighborhoods of $\Gamma$ by $\left[{ }^{7}\right]$, and that $\operatorname{grad} K_{D}$ maps into $\bigcap_{r \in \mathbf{R}} H_{p}^{r}(\{|x| \geq R\})^{n}$ when $R$ is large enough. For $\operatorname{grad} R_{D}$ div, the property is straightforward when $s \geq 0$; to include lower $s$, one uses that the operator is selfadjoint (being an orthogonal projection in $L_{2}$ ), and that $H_{p}^{(s, s / 2)}(\bar{\Omega})^{*}=$ $H_{p^{\prime}}^{(-s,-s / 2)}(\bar{\Omega})^{*}$ for $\left.s \in\right] \frac{1}{p}-1, \frac{1}{p}\left[\right.$, here $\frac{1}{p}+\frac{1}{p^{\prime}}=1$.)

For the other boundary conditions (the cases $k=0,2,3,4$ ), the above analysis can be carried through in suitably modified versions.

The application of the linear result to solve the nonlinear problem goes mechanically as in the bounded case in $[7]$, so we arrive at the result:

Theorem 4.3 Theorem 2.1 generalizes to the case of an exterior smooth domain.
(When $p>n /(n-1), \operatorname{grad} q$ is only unique up to addition of functions $c(t) \operatorname{grad}|x|^{2-n}$ for $n>2$, resp. $c(t) \operatorname{grad} \log |x|$ for $n=2$.)

## Acknowledgments

The author is grateful to B. Fuglede and to M. Yamazaki for useful conversations.

## References

1. M. Brelot: Éléments de la théorie classique du potentiel. Les cours de Sorbonne, 3e cycle, 3e Édition (Centre de Documentation Universitaire, Paris, 1965).
2. Y. Giga: The nonstationary Navier-Stokes system with some first order boundary conditions, Proc. Jap. Acad. 58, 101 (1982).
3. Y. Giga and T. Miyakawa: Solutions in $L_{r}$ of the Navier-Stokes initial value problem, Arch. Rat. Mech. Anal. 189, 267 (1985).
4. G. Grubb: Functional Calculus of Pseudo-Differential Boundary Problems, Progress in Math. Vol. 65 (Birkhäuser, Boston, 1986).
5. G. Grubb: Pseudo-differential boundary problems in $L_{p}$ spaces, Comm. P. D. E. 15, 289 (1990).
6. G. Grubb: Parameter-elliptic and parabolic pseudodifferential boundary problems in global $L_{p}$ Sobolev spaces, Math. Zeitschr. 218, 43 (1995).
7. G. Grubb: Nonhomogeneous time-dependent Navier-Stokes problems in $L_{p}$ Sobolev spaces, Differential and Integral Equations 8, 1013 (1995).
8. G. Grubb: Initial value problems for the Navier-Stokes equations with Neumann conditions, in The Navier-Stokes Equations II - Theory and Numerical Methods, Proceedings Oberwolfach 1991, p. 262 (Springer Lecture Note no. 1530, Heidelberg, 1992) editors J. G. Heywood et al.
9. G. Grubb and N. J. Kokholm: A global calculus of parameter-dependent pseudodifferential boundary problems in $L_{p}$ Sobolev spaces Acta Mathematica 171, 165 (1993).
10. G. Grubb and R. Seeley: Weakly parametric pseudodifferential operators and Atiyah-Patodi-Singer boundary problems, Inventiones Math. 121, 481 (1995).
11. G. Grubb and V. A. Solonnikov: Reduction of basic initial-boundary value problems for the Stokes equation to initial-boundary value problems for systems of pseudodifferential equations Zapiski Nauchn. Sem. L.O.M.I. 163, 37 (1987) = J. Soviet Math. 49, 1140 (1990).
12. G. Grubb and V. A. Solonnikov: Solution of parabolic pseudo-differential initial-boundary value problems, J. Diff. Equ. 87, 256 (1990).
13. G. Grubb and V. A. Solonnikov: Reduction of basic initial-boundary value problems for the Navier-Stokes equations to nonlinear parabolic systems of pseudodifferential equations, Zap. Nauchn. Sem. L.O.M.I. 171, 36 (1989) $=$ J. Soviet Math. 56, 2300 (1991).
14. G. Grubb and V. A. Solonnikov: Boundary value problems for the nonstationary Navier-Stokes equations treated by pseudo-differential methods, Math. Scand. 1991, 217 (69).
15. O. A. Ladyzhenskaya, V. A. Solonnikov and N. N. Uralceva: Linear and Quasi-linear Parabolic Equations, AMS Translation Math. Monographs 23 (AMS, Providence, Rhode Island 1968).
16. C. G. Simader and H. Sohr: The Weak Dirichlet Problem for $\Delta$ in $L^{q}$ in Bounded and Exterior Domains (Pitman Research Notes, London), to appear.
17. V. A. Solonnikov: Estimates of the solutions of a nonstationary linearized system of Navier-Stokes equations, Trudy Mat. Inst. Steklov 70, 213 (1964) $=A M S$ Translations 75, 1 (1968)
18. V. A. Solonnikov: Estimates for solutions of nonstationary Navier-Stokes systems Zap. Nauchn. Sem. LOMI 38, 153 (1973) = J. Soviet Math. 8, 467 (1977)
19. V. A. Solonnikov: Estimates of solutions of an initial- and boundaryvalue problem for the linear nonstationary Navier-Stokes system Zap. Nauchn. Sem. LOMI 59, $172(1976)=$ J. Soviet Math. 10, 336 (1978)
20. W. von Wahl: The Equations of Navier-Stokes and Abstract Parabolic Equatios (Vieweg und Sohn, Braunschweig, 1985).
