QUALITATIVE BEHAVIOUR OF SOLUTIONS
FOR THE TWO-PHASE NAVIER-STOKES EQUATIONS
WITH SURFACE TENSION

MATTHIAS KÖHNE, JAN PRÜSS, AND MATHIAS WILKE

In memory of Professor Tetsuro Miyakawa †

Abstract. The two-phase free boundary value problem for the isothermal
Navier-Stokes system is studied for general bounded geometries in absence of
phase transitions, external forces and boundary contacts. It is shown that the
problem is well-posed in an $L^p$-setting, and that it generates a local semiflow on
the induced phase manifold. If the phases are connected, the set of equilibria of
the system forms a $(n+1)$-dimensional manifold, each equilibrium is stable, and
it is shown that global solutions which do not develop singularities converge
to an equilibrium as time goes to infinity. The latter is proved by means of
the energy functional combined with the generalized principle of linearized
stability.

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ity, compactness, generalized principle of linearized stability, convergence.

1. Introduction

In this paper we consider a free boundary problem that describes the motion
of two isothermal, viscous, incompressible Newtonian fluids in $\mathbb{R}^3$. The fluids are
separated by an interface that is unknown and has to be determined as part of the
problem.

More precisely, we consider two fluids that fill a region $\Omega \subset \mathbb{R}^3$. Let $\Gamma_0 \subset \Omega$ be a
given surface which bounds the region $\Omega_1(0)$ occupied by a viscous incompressible
fluid, fluid$_1$, the dispersed phase and let $\Omega_2(0)$ be the complement of the closure of
$\Omega_1(0)$ in $\Omega$, corresponding to the region occupied by a second incompressible viscous
fluid, fluid$_2$, the continuous phase. Note that the dispersed phase is assumed not to
be in contact with the boundary $\partial \Omega$ of $\Omega$. We assume that the two fluids are immiscible, and that no phase transitions occur. The velocity of the fluids is
denoted by $u(t, x)$, and the pressure field by $\pi(t, x)$.

Let $\Gamma(t)$ denote the position of $\Gamma_0$ at time $t$. Thus, $\Gamma(t)$ is a sharp interface which
separates the fluids occupying the regions $\Omega_1(t)$ and $\Omega_2(t)$, respectively. We denote
the normal field on $\Gamma(t)$, pointing from $\Omega_1(t)$ into $\Omega_2(t)$, by $n_\Gamma(t, \cdot)$. Moreover,$V_\Gamma(t, \cdot)$ and $H_\Gamma(t, \cdot)$ mean the normal velocity and the curvature of $\Gamma(t)$ with respect
to $n_\Gamma(t, \cdot)$, respectively. Here the curvature $H_\Gamma := -\text{div}_\Gamma V_\Gamma$ is negative when $\Omega_1(t)$
is convex in a neighborhood of $x \in \Gamma(t)$, in particular the curvature of a sphere
\[ S_R(x_0) = -(n-1)/R. \] The motion of the fluids is governed by the following system of equations for \( i = 1, 2. \)

\[ \rho_i \left( \partial_t u + (u \nabla)u \right) - \mu_i \Delta u + \nabla \pi = 0 \quad \text{in} \quad \Omega_i(t), \]
\[ \text{div} u = 0 \quad \text{in} \quad \Omega_i(t), \]
\[ -[S(u, \pi)\nu_T] = \sigma H \nu_T \quad \text{on} \quad \Gamma(t), \]
\[ [u] = 0 \quad \text{on} \quad \Gamma(t), \]
\[ V_T = (u|\nu_T) \quad \text{on} \quad \Gamma(t), \]
\[ u(0) = u_0 \quad \text{in} \quad \Omega \setminus \Gamma_0, \]
\[ \Gamma(0) = \Gamma_0. \]

(1.1)

Here, \( S \) is the stress tensor defined by
\[ S = S(u, \pi) = \mu_i (\nabla u + [\nabla u]^T) - \pi I = 2\mu_i E - \pi I \quad \text{in} \quad \Omega_i(t), \]
and
\[ [\phi](t, x) = \lim_{h \to 0^+} \left( \phi(t, x + h\nu_T(x)) - \phi(t, x - h\nu_T(x)) \right), \quad x \in \Gamma(t) \]
denotes the jump of the quantity \( \phi \), defined on the respective domains \( \Omega_i(t) \), across the interface \( \Gamma(t) \).

Given are the initial position \( \Gamma_0 \) of the interface and the initial velocity \( u_0 : \Omega \setminus \Gamma_0 \to \mathbb{R}^3 \). The unknowns are the velocity field \( u(t, \cdot) : \Omega \setminus \Gamma(t) \to \mathbb{R}^3 \), the pressure field \( \pi(t, \cdot) : \Omega \setminus \Gamma(t) \to \mathbb{R} \), and the free boundary \( \Gamma(t) \).

The constants \( \rho_i > 0 \) and \( \mu_i > 0 \) denote the densities and the viscosities of the respective fluids, and the constant \( \sigma > 0 \) stands for the surface tension. In the sequel we drop the index \( i \) since there is no danger of confusion; however, we keep in mind that \( \mu \) and \( \rho \) have jumps across the interface, in general.

System (1.1) comprises the two-phase Navier-Stokes equations with surface tension. The corresponding one-phase problem is obtained by setting \( \rho_2 = \mu_2 = 0 \) and discarding \( \Omega_2 \). Here we concentrate the discussion on the two-phase problem.

There are several papers in the literature dealing with problem (1.1); cf. [5, 6, 7, 8, 9, 25, 26, 27]. All of them employ Lagrangian coordinates to obtain local well-posedness. This way it seems difficult to establish smoothing of the unknown interface, and this method is hardly useful in case phase transitions have to be taken into account. Here we employ a different approach, namely the Direct Mapping Method via a Hanzawa transform, which has been quite efficient in the study of Stefan problems, i.e. phase transitions involving temperature, only.

In a recent paper [19] we have shown that problem (1.1) is locally well-posed in an \( L_p \)-setting provided \( \Omega = \mathbb{R}^n \) and the initial interface \( \Gamma_0 \) is sufficiently close to a flat configuration. In addition, the interface as well as the solution are proved to become instantaneously real analytic. This result is based on a careful analysis of the underlying linear problem. Building on the latter results we show in this paper local well-posedness for arbitrary bounded geometries as described above. This induces a local semiflow on a well-defined nonlinear phase manifold.

It is known that the set \( E \) of equilibria of the system are zero velocities, constant pressures in the components of the phases and the dispersed phase is a union of disjoint balls. Concentrating on the case of connected phases, we prove that
such that $G \eta > 0$ there is a real analytic closed hypersurface $\Sigma$, in the sense that if $\eta > 0$ is small enough, then $\Sigma$ bounds a domain $G_1$ with $\overline{G_1} \subset \Omega$, and we set $G_2 = \Omega \setminus G_1$.

It is well known that such a hypersurface $\Sigma$ admits a tubular neighbourhood, which means that there is $a > 0$ such that the map

$$
\Lambda : \Sigma \times (-a, a) \to \mathbb{R}^n, \quad \Lambda(p, r) := p + r\nu_\Sigma(p), \quad p \in \Sigma, \ |r| < a
$$

is a diffeomorphism from $\Sigma \times (-a, a)$ onto $\mathcal{R}(\Lambda)$. The inverse

$$
\Lambda^{-1} : \mathcal{R}(\Lambda) \to \Sigma \times (-a, a)
$$

of this map is conveniently decomposed as

$$
\Lambda^{-1}(x) = (\Pi(x), d(x)), \quad x \in \mathcal{R}(\Lambda).
$$

Here $\Pi(x)$ means the orthogonal projection of $x$ to $\Sigma$ and $d(x)$ the signed distance from $x$ to $\Sigma$; so $|d(x)| = \text{dist}(x, \Sigma)$ and $d(x) < 0$ iff $x \in G_1$. In particular we have

$$
\mathcal{R}(\Lambda) = \{x \in \mathbb{R}^n : \text{dist}(x, \Sigma) < a\}.
$$

Note that on the one hand, $a$ is determined by the curvatures of $\Sigma$, i.e. we must have

$$
0 < a < \min\{1/|\kappa_j(p)| : j = 1, \ldots, n-1, \ p \in \Sigma\},
$$

where $\kappa_j(p)$ mean the principal curvatures of $\Sigma$ at $p \in \Sigma$. But on the other hand, $a$ is also connected to the topology of $\Sigma$, which can be expressed as follows. Since $\Sigma$ is a compact manifold of dimension $n-1$ it satisfies the ball condition, which means that there is a radius $r_\Sigma > 0$ such that for each point $p \in \Sigma$ there are $x_j \in G_j$, $j = 1, 2$, such that $B_{r_\Sigma}(x_j) \subset G_j$, and $B_{r_\Sigma}(x_j) \cap \Sigma = \{p\}$. Choosing $r_\Sigma$ maximal, we then must also have $a < r_\Sigma$. 

2. Transformation to a Fixed Domain

Let $\Omega \subset \mathbb{R}^n$ be a bounded domain with boundary $\partial \Omega$ of class $C^2$, and suppose $\Gamma \subset \Omega$ is a hypersurface of class $C^2$, i.e. a $C^2$-manifold which is the boundary of a bounded domain $\Omega_1 \subset \Omega$; we then set $\Omega_2 = \Omega \setminus \Omega_1$. Note that $\Omega_2$ is connected, but $\Omega_1$ maybe disconnected, however, it consists of finitely many components only, since $\partial \Omega_1 = \Gamma$ by assumption is a manifold, at least of class $C^2$. Recall that the second order bundle of $\Gamma$ is given by

$$
N^2\Gamma := \{(p, \nu_\Gamma(p), \nabla_\Gamma \nu_\Gamma(p)) : p \in \Gamma\}.
$$

Here $\nabla_\Gamma$ denotes the surface gradient on $\Gamma$. Recall also the Hausdorff distance $d_H$ between the two closed subsets $A, B \subset \mathbb{R}^n$, defined by

$$
d_H(A, B) := \max\{\sup_{a \in A} \text{dist}(a, B), \sup_{b \in B} \text{dist}(b, A)\}.
$$

Then we may approximate $\Gamma$ by a real analytic hypersurface $\Sigma$, in the sense that the Hausdorff distance of the second order bundles of $\Gamma$ and $\Sigma$ is as small as we want. More precisely, for each $\eta > 0$ there is a real analytic closed hypersurface $\Sigma$ such that $d_H(N^2\Sigma, N^2\Gamma) \leq \eta$. If $\eta > 0$ is small enough, then $\Sigma$ bounds a domain $G_1$ with $\overline{G_1} \subset \Omega$, and we set $G_2 = \Omega \setminus G_1$.

It is well known that such a hypersurface $\Sigma$ admits a tubular neighbourhood, which means that there is $a > 0$ such that the map

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where $\kappa_j(p)$ mean the principal curvatures of $\Sigma$ at $p \in \Sigma$. But on the other hand, $a$ is also connected to the topology of $\Sigma$, which can be expressed as follows. Since $\Sigma$ is a compact manifold of dimension $n-1$ it satisfies the ball condition, which means that there is a radius $r_\Sigma > 0$ such that for each point $p \in \Sigma$ there are $x_j \in G_j$, $j = 1, 2$, such that $B_{r_\Sigma}(x_j) \subset G_j$, and $B_{r_\Sigma}(x_j) \cap \Sigma = \{p\}$. Choosing $r_\Sigma$ maximal, we then must also have $a < r_\Sigma$. 

The energy of the system serves as a strict Ljapunov functional, hence the limit sets of the solutions are contained in the set of equilibria $\mathcal{E}$. Combining these results we show that any solution which does not develop singularities converges to an equilibrium in the topology of the phase manifold.
Setting $\Gamma(0) = \Gamma_0$, we may use the map $\Lambda$ to parametrize the unknown free boundary $\Gamma(t)$ over $\Sigma$ by means of a height function $h$ via

$$\Gamma(t) = R(p \mapsto p + h(t, p)\nu_\Sigma(p), \quad p \in \Sigma), \quad t \geq 0,$$

for small $t \geq 0$, at least. Extend this diffeomorphism to all of $\bar{\Omega}$ by means of

$$\Theta_h(t, x) = x + \chi(d(x)/a)h(t, \Pi(x))\nu_\Sigma(\Pi(x)) =: x + \theta_h(t, x).$$

Here $\chi$ denotes a suitable cut-off function; more precisely, $\chi \in D(\mathbb{R})$, $0 \leq \chi \leq 1$, $\chi(r) = 1$ for $|r| < 1/3$, and $\chi(r) = 0$ for $|r| > 2/3$. This way $\Omega \backslash \Gamma(t)$ is transformed to the fixed domain $\Omega \backslash \Sigma$. Note that $\Theta_h(t, x) = x$ for $|d(x)| > 2a/3$, and

$$\Theta_h^{-1}(t, x) = x - h(t, \Pi(x))\nu_\Sigma(\Pi(x)) \quad \text{for} \quad |d(x)| < a/3,$$

in particular,

$$\Theta_h^{-1}(t, x) = x - h(t, x)\nu_\Sigma(x) \quad \text{for} \quad x \in \Sigma.$$

Now we define the transformed quantities

$$\bar{u}(t, x) = u(t, \Theta_h(t, x)), \quad \bar{\pi}(t, x) = \pi(t, \Theta_h(t, x)), \quad t > 0, \quad x \in \Omega \backslash \Sigma,$$

the pull backs of $u$ and $\pi$. This gives the following problem for $\bar{u}, \bar{\pi}, h$.

$$\rho \partial_t \bar{u} - \mu A(h)\bar{u} + G(h)\bar{\pi} = R(\bar{u}, h) \quad \text{in} \quad \Omega \backslash \Sigma, \quad (G(h)|\bar{u}) = 0 \quad \text{in} \quad \Omega \backslash \Sigma,$$

$$-[\mu([G(h)\bar{u}] + [G(h)\bar{u}]^T) - \bar{\pi}]\nu_\Gamma(h) = \sigma H_\Gamma(h)\nu_\Gamma(h) \quad \text{on} \quad \Sigma, \quad [\bar{u}] = 0 \quad \text{on} \quad \Sigma,$$

$$\bar{u} = 0 \quad \text{on} \partial \bar{\Omega}, \quad \partial_t h - (\bar{u})\nu_\Sigma = -((\bar{u})|\alpha(h)) \quad \text{on} \quad \Sigma,$$

$$\bar{u}(0) = \bar{u}_0, \quad \text{in} \quad \Omega \backslash \Sigma, \quad h(0) = h_0 \quad \text{on} \quad \Sigma.$$

Here $A(h), G(h)$ and $H_\Gamma(h)$ denote the transformed Laplacian, gradient and curvature, respectively. More precisely, we have

$$\Theta_h' = I + \theta_h', \quad \Theta_h^{-1} = I - [I + \theta_h']^{-1}\theta_h'$$

and

$$[\nabla \pi] \circ \Theta_h = G(h)\bar{\pi} = ([\Theta_h^{-1}]^T \circ \Theta_h)[\nabla \bar{\pi}] = \nabla \bar{\pi} - \theta_h' [I + \theta_h']^{-T} \nabla \bar{\pi} =: (I - M_1(\bar{h}))\nabla \bar{\pi}$$

$$[\text{div } u] \circ \Theta_h = (G(h)|\bar{u}) = ((\Theta_h^{-1})^T \circ \Theta_h) \nabla |\bar{u}| = (\nabla |\bar{u}|) - (\theta_h' [I + \theta_h']^{-T} \nabla |\bar{u}|) = ((I - M_1(\bar{h}))\nabla |\bar{u}|$$

and

$$[\Delta u] \circ \Theta_h = A(h)\bar{u} = ([\Theta_h^{-1}]^T (\Theta_h^{-1})^T \circ \Theta_h) : \nabla^2 \bar{u} + ([\Delta \Theta_h^{-1}] \circ \Theta_h) \nabla \bar{u} = \Delta \bar{u} - M_2(h) \nabla \bar{u}$$
Note that
\[ P \text{ with } H \]
the curvature!

With the curvature tensor \( h \), hence
\[ \text{linearization is given by} \]
\[ \text{Employing this notation, we have} \]
\[ \text{and} \]
\[ \text{hence} \]
\[ R(\bar{u}, h) = -\rho(\bar{u} \cdot \mathcal{G}(h)\bar{u}) + M_3(h)\nabla \bar{u}. \]

With the curvature tensor \( L_\Sigma \) and the surface gradient \( \nabla \Sigma \) we have
\[ \nu_T(h) = \beta(h)(\nu_\Sigma - \alpha(h)), \quad \alpha(h) = M_0(h)\nabla h, \]
\[ M_0(h) = (I - hL_\Sigma)^{-1}, \quad \beta(h) = (1 + |\alpha(h)|^2)^{-1/2}, \]
and
\[ V = (\partial_t \Theta | \nu_T) = \partial_t h(\nu_T | \nu_\Sigma) = \beta(h)\partial_t h. \]

Employing this notation, we have
\[ \theta_\Sigma'(t, x) = \nu_\Sigma(\Pi(x)) \otimes M_0(\varphi(x))\nabla h(t, \Pi(x)) \]
\[ - h(t, \Pi(x))L_\Sigma(\Pi(x))M_0(\varphi(x))\mathcal{P}_\Sigma \quad \text{for} \quad |\varphi(x)| < a/3, \]
\[ \theta_\Sigma'(t, x) = 0 \quad \text{for} \quad |\varphi(x)| > 2a/3, \]
and
\[ \theta'_\Sigma(t, x) = \frac{1}{a} \theta'(d(x)/a)h(t, \Pi(x))\nu_\Sigma(\Pi(x)) \otimes \nu_\Sigma(\Pi(x)) \]
\[ + \theta'(d(x)/a)\nu_\Sigma(\Pi(x)) \otimes M_0(\varphi(x))\nabla h(t, \Pi(x)) \]
\[ - \theta'(d(x)/a)h(t, \Pi(x))L_\Sigma(\Pi(x))M_0(\varphi(x))\mathcal{P}_\Sigma \]
\[ \quad \text{for} \quad a/3 < |\varphi(x)| < 2a/3, \]
where \( \mathcal{P}_\Sigma = I - \nu_\Sigma \otimes \nu_\Sigma \) denotes the projection onto the tangent space of \( \Sigma \). Thus, \([I + \theta_\Sigma] \) is boundedly invertible, if \( |h|_\infty \) and \( |\nabla \Sigma h|_\infty \) are sufficiently small. The curvature \( H_\Gamma(h) \) becomes
\[ H_\Gamma(h) = \beta(h)([\text{tr}[M_0(h)(L_\Sigma + \nabla \Sigma \alpha(h))]) - \beta^2(h)(M_0(h)\alpha(h)||\nabla \Sigma \alpha(h)||\alpha(h))), \]
a differential expression involving second order derivatives of \( h \) only linearly. Its linearization is given by
\[ H_\Gamma^T(0) = \text{tr} \, L_\Sigma^2 + \Delta_\Sigma. \]

Here \( \Delta_\Sigma \) denotes the Laplace-Beltrami operator on \( \Sigma \).

It is convenient to decompose the stress boundary condition into tangential and normal parts. Multiplying the stress interface condition with \( \nu_\Sigma/\beta \) we obtain
\[ [\bar{\sigma}] - \sigma H_\Gamma(h) = ([\mu(\mathcal{G}(h)\bar{u}) + [\mathcal{G}(h)\bar{u}^T])(\nu_\Sigma - M_0(h)\nabla h|\nu_\Sigma), \]
for the normal part of the stress boundary condition, and
\[ -\mathcal{P}_\Sigma\mu([\mathcal{G}(h)\bar{u}) + [\mathcal{G}(h)\bar{u}^T])(\nu_\Sigma - M_0(h)\nabla h) \]
\[ = ([\mu(\mathcal{G}(h)\bar{u}) + [\mathcal{G}(h)\bar{u}^T])(\nu_\Sigma - M_0(h)\nabla h|\nu_\Sigma) M_0(h)\nabla h, \]
for the tangential part. Note that the latter neither contains the pressure jump nor the curvature!
We rewrite this problem in quasilinear form, dropping the bars and collecting its principal linear part on the left hand side.

\[
\rho \partial_t u - \mu \Delta u + \nabla \pi = F(h, u) \nabla u + M_4(h) : \nabla^2 u + M_1(h) \nabla \pi \quad \text{in } \Omega \setminus \Sigma,
\]
\[
\text{div } u = M_1(h) : \nabla u \quad \text{in } \Omega \setminus \Sigma,
\]
\[
P_\Sigma \left[ - \mu (\nabla u + \nabla u^T) \right] \nu_\Sigma = G_r(h) \nabla u, \quad \text{(2.2)}
\]
\[
- \big( - \mu (\nabla u + \nabla u^T) \big) \nu_\Sigma \nu_\Sigma + \left[ \nu \right] - \sigma H_1^1(0) h = G_\nu(h) \nabla u + \Gamma_g(h) \quad \text{on } \Sigma,
\]
\[
[u] = 0 \quad \text{on } \Sigma,
\]
\[
u = 0 \quad \text{on } \partial \Omega,
\]
\[
\partial_t h - (u|\nu_\Sigma) = (M_0(h) \nabla_\Sigma h|u) \quad \text{on } \Sigma,
\]
\[
u(0) = u_0 \quad \text{in } \Omega \setminus \Sigma, \quad h(0) = h_0 \quad \text{on } \Sigma.
\]

The right-hand sides in this problem are either lower order terms or are of the same order appearing on the left, but carrying a factor \( h \) or \( \nabla_\Sigma h \), which are small by construction. In fact, since \( \Gamma_0 \) is approximated by \( \Sigma \) in the second order bundle we have smallness of \( h_0, \nabla_\Sigma h_0 \), and even of \( \nabla_\Sigma^2 h_0 \), uniformly on \( \Sigma \). All terms on the right-hand side are at least quadratic. More precisely, besides the \( M_j(h) \) which have been introduced before, the nonlinearities have the following form:

\[
F(h, u) \nabla u = - (u|\nabla u) + [M_1(h) + M_2(h) + M_3(h)] \nabla u,
\]
\[
G_\nu(h) \nabla u = \left( \mu |[\nabla u] + [\nabla u]^T| \right) [M_0(h) \nabla_\Sigma h|\nu_\Sigma) 
\]
\[
- \big( \mu \left( [M_1(h) \nabla u] + [M_1(h) \nabla u]^T \right) \big) (\nu_\Sigma - M_0(h) \nabla_\Sigma h|\nu_\Sigma) ,
\]
\[
G_r(h) \nabla u = \big( \mu \left( [(I - M_1(h)) \nabla u] + [(I - M_1(h)) \nabla u]^T \right) \big) (\nu_\Sigma - M_0(h) \nabla_\Sigma h|\nu_\Sigma) .
\]
\[
\cdot M_0(h) \nabla_\Sigma h 
\]
\[
- P_\Sigma \mu \left( [(I - M_1(h)) \nabla u] + [(I - M_1(h)) \nabla u]^T \right) [M_0(h) \nabla_\Sigma h 
\]
\[
- P_\Sigma \mu \left( [M_1(h) \nabla u] + [M_1(h) \nabla u]^T \right) [\nu_\Sigma, 
\]
\[
G_\nu(h) = \sigma (H_1^1(h) - H_1^1(0) h).
\]

The idea of our approach can be described as follows. We consider the transformed problem \((2.2)\). Based on maximal \( L_p \)-regularity of the linear problem given by the left hand side of \((2.2)\), we employ the contraction mapping principle to obtain local well-posedness of the nonlinear problem. The solutions of the transformed problem will belong to the following class:

\[
u \in H_1^1(J; L_p(\Omega^n)) \cap L_p(J; H_1^2(\Omega \setminus \Sigma)^n), \quad \pi \in L_p(J; \dot{H}_1^2(\Omega \setminus \Sigma)),
\]
\[
[\pi] \in W^{1/2-1/2p}_p(J; L_p(\Sigma)) \cap L_p(J; W^{1-1/p}_p(\Sigma)), 
\]
\[
h \in W^{2-1/2p}_p(J; L_p(\Sigma)) \cap H^1_p(J; W^{2-1/p}_p(\Sigma)) \cap L_p(J; W^{3-1/p}_p(\Sigma)).
\]

This program will be carried out in the next sections.
3. The Linearized Problem

We consider now the inhomogeneous linear problem

\[
\begin{align*}
\rho \partial_t u - \mu \Delta u + \nabla \pi &= \rho f & \text{in } \Omega \setminus \Sigma, \\
\text{div } u &= f_d & \text{in } \Omega \setminus \Sigma, \\
[ -\mu(\nabla u + [\nabla u]^T) + I \pi ]\nu_{\Sigma} - \sigma(\Delta_{\Sigma} h)\nu_{\Sigma} &= g & \text{on } \Sigma, \\
[u] &= u_{\Sigma} & \text{on } \Sigma, \\
\langle u \rangle &= u_b & \text{on } \partial \Omega,
\end{align*}
\]

(3.1)

on a finite time-interval \( J = [0, a] \). We choose the same regularity classes for \( u \) and \( \pi \) as before, i.e.

\[
u_{\Sigma} \in W_p^{1-1/2p}(J; L_p(\Sigma)^n) \cap L_p(J; H_p^2(\Sigma)^n),
\]

and

\[
\pi \in Z_{\Sigma} := L_p(J; \hat{H}_p^3(\Sigma)).
\]

Then

\[
u_{\Sigma} \in W_p^{1-1/2p}(J; L_p(\Sigma)^n) \cap L_p(J; W_p^{2-1/p}(\Sigma)^n),
\]

and

\[
u_b \in W_p^{1-1/2p}(J; L_p(\partial \Omega)^n) \cap L_p(J; W_p^{2-1/p}(\partial \Omega)^n).
\]

Therefore the equation for the height function \( h \) lives in the trace space for the components of \( u \), i.e.

\[
g_h \in Y_u^0 := W_p^{1-1/2p}(J; L_p(\Sigma)) \cap L_p(J; W_p^{2-1/p}(\Sigma)),
\]

hence the natural space for \( h \) is given by

\[
h \in Z_h := W_p^{2-1/2p}(J; L_p(\Sigma)) \cap H_p^1(J; W_p^{2-1/p}(\Sigma)) \cap L_p(J; W_p^{3-1/p}(\Sigma)).
\]

Here the last space comes from the curvature term in the stress boundary condition, which induces an additional order in spatial regularity. Assuming that \( g \) belongs to the trace space of \( \nabla u \), i.e.

\[
g \in Y_u^1 := W_p^{1-2/2p}(J; L_p(\Sigma)^n) \cap L_p(J; W_p^{1-1/p}(\Sigma)^n),
\]

we have the additional regularity \( \| \pi \| \in Y_u^1 \) for the pressure jump across the interface \( \Sigma \). The function \( b \in Y_u^0 \) is given; we will choose \( b \) appropriately in Section 4.

There is another hidden regularity which comes from the divergence equation. To identify it, let \( \phi \in \hat{H}_p^1(\Omega) \). An integration by parts yields

\[
\begin{align*}
(u|\nabla \phi)_\Omega &= -(\text{div } u|\phi)_\Omega + (u \cdot \nu_{\partial \Omega}|\phi)_{\partial \Omega} - ([u \cdot \nu_{\Sigma}]|\phi)_{\Sigma} \\
&= -(f_d|\phi)_\Omega + (u_b \cdot \nu_{\partial \Omega}|\phi)_{\partial \Omega} - (u_{\Sigma} \cdot \nu_{\Sigma}|\phi)_{\Sigma}.
\end{align*}
\]

Set \( \hat{H}_p^{-1}(\Omega) = (\hat{H}_p^1(\Omega))^* \) and define the functional \( (f_d, u_b \cdot \nu_{\partial \Omega}, u_{\Sigma} \cdot \nu_{\Sigma}) \in \hat{H}_p^{-1}(\Omega) \) by means of

\[
\langle (f_d, u_b \cdot \nu_{\partial \Omega}, u_{\Sigma} \cdot \nu_{\Sigma})|\phi \rangle := -(f_d|\phi)_\Omega + (u_b \cdot \nu_{\partial \Omega}|\phi)_{\partial \Omega} - (u_{\Sigma} \cdot \nu_{\Sigma}|\phi)_{\Sigma}.
\]

Then we have

\[
\langle (f_d, u_b \cdot \nu_{\partial \Omega}, u_{\Sigma} \cdot \nu_{\Sigma})|\phi \rangle = (u|\nabla \phi)_\Omega.
\]
Since \( u \in H^1_p(J; L_p(\Omega)^n) \) this implies \((f_d, u_0 \cdot \nu_{\partial \Omega}, u_\Sigma \cdot \nu_\Sigma) \in H^1_p(J; \tilde{H}^{-1}_p(\Omega)) \). Observe that this condition contains the compatibility condition
\[
\int_{\Omega} f_d \, dx = \int_{\partial \Omega} u_0 \cdot \nu_{\partial \Omega} \, d\Omega - \int_{\Sigma} u_\Sigma \cdot \nu_\Sigma \, d\Sigma,
\]
which appears choosing \( \phi = 1 \).

In the particular case \( f_d = 0 \) we have \((f_d, u_0 \cdot \nu_{\partial \Omega}, u_\Sigma \cdot \nu_\Sigma) \in H^1_p(J; \tilde{H}^{-1}_p(\Omega)) \) if and only if \( u_0 \cdot \nu_{\partial \Omega} \in H^1_p(J; W^{-1/p}_p(\partial \Omega)) \) and \( u_\Sigma \cdot \nu_\Sigma \in H^1_p(J; W^{-1/p}_p(\Sigma)) \).

The main theorem of this section states that problem \( 3.1 \) admits maximal regularity, in particular, it defines an isomorphism between the solution space and the space of data.

**Theorem 3.1.** Let \( p > n + 2 \), \( \Omega \subset \mathbb{R}^n \) a bounded domain with \( \partial \Omega \subset C^3 \), \( \Sigma \subset \Omega \) a closed hypersurface of class \( C^3 \) and \( \rho_j \), \( \mu_j \), \( \sigma \) be positive constants, \( j = 1, 2 \); set \( J = [0, a] \), and suppose
\[
b \in W^{-1/2p}_p(J; L_p(\Sigma))^n \cap L_p(J; W^{2-1/p}_p(\Sigma))^n.
\]
Then the two-phase Stokes problem \( 3.1 \) admits a unique solution \((u, \pi, h)\) with regularity
\[
\begin{align*}
&u \in H^1(J; L_p(\Omega)^n) \cap L_p(J; H^{2}_p(\Omega \setminus \Sigma)^n), \quad \pi \in L_p(J; H^{1}_p(\Sigma)), \\
&[\pi] \in W^{1/2-1/2p}_p(J; L_p(\Sigma)) \cap L_p(J; W^{1-1/2p}_p(\Sigma)), \\
h \in W^{2-1/2p}_p(J; L_p(\Sigma))^n \cap H^{1}_p(J; W^{2-1/p}_p(\Sigma))^n \cap L_p(J; W^{3-1/p}_p(\Sigma)),
\end{align*}
\]
if and only if the data \((u_0, h_0, u_0, u_\Sigma, f, f_d, g, g_h)\) satisfy the following regularity and compatibility conditions:
(a) \( f \in L_p(J \times \Omega)^n \), \( u_0 \in W^{2-2/p}_p(\Omega \setminus \Sigma)^n \), and \( u_0|_{\partial \Omega} = u_{0|_{\partial \Omega}} \); (b) \( f_d \in L_p(J; H^1_p(\Omega \setminus \Sigma)) \), and \( \text{div } u_0 = f_d|_{t=0} \); (c) \( u_\Sigma \in W^{1-1/2p}_p(J; L_p(\Sigma))^n \cap L_p(J; W^{2-1/p}_p(\partial \Omega)^n) \); (d) \( u_\Sigma \in W^{1-1/2p}_p(J; L_p(\Sigma))^n \cap L_p(J; W^{2-1/p}_p(\Sigma))^n) \); (e) \( f_d, u_0 \cdot \nu_{\partial \Omega}, u_\Sigma \cdot \nu_\Sigma \in H^{1}_p(J; \tilde{H}^{-1}_p(\Omega)) \); (f) \( g \in W^{1/2-1/2p}_p(J; L_p(\Sigma))^n \cap L_p(J; W^{2-1/p}_p(\Sigma))^n) \); (g) \( [u_0] = u_{0|_{t=0}}, \) and \( P_\Sigma[\mu(\nabla u_0 + (\nabla u_0)^T)] = P_\Sigma g_{\mid_{t=0}} \); (h) \( h_0 \in W^{3-2/p}_p(\Sigma), \) and \( g_h \in W^{1-1/2p}_p(J; L_p(\Sigma))^n \cap L_p(J; W^{2-1/2p}_p(\Sigma)).

The solution map \((u_0, h_0, u_0, u_\Sigma, f, f_d, g, g_h, b) \mapsto (u, \pi, [\pi], h)\) is continuous between the corresponding spaces.

The proof will be carried out in the following subsections.

In general the pressure \( \pi \) has no more regularity as stated in Theorem 3.1. However, there are situations where \( \pi \) enjoys extra time-regularity, as stated in the following

**Corollary 3.2.** Assume in addition to the hypotheses of Theorem 3.1 that
\[
\begin{align*}
&u_0 = h_0 = f_d = 0, \quad \text{div } f = 0 \text{ in } \Omega \setminus \Sigma, \\
&u_\Sigma \cdot \nu_\Sigma = 0 \text{ on } \partial \Omega, \quad u_\Sigma \cdot \nu_\Sigma = 0 \text{ on } \Sigma,
\end{align*}
\]
and
\[
[[f|\nu_\Sigma]] = 0 \text{ on } \Sigma, \quad (f|\nu_{\partial \Omega}) = 0 \text{ on } \partial \Omega.
\]
Then \( \pi \in H^\alpha_p(J; L_p(\Omega)) \), for each \( \alpha \in (0, 1/2 - 1/2p) \).
The function \( \pi \) is determined by the data, and we define \( \pi \) as the solution of the parabolic problem

\[
\begin{align*}
\Delta \phi &= \rho q \quad \text{in } \Omega \setminus \Sigma, \\
[\phi] &= 0 \quad \text{on } \Sigma, \\
[\rho^{-1} \partial_{\nu} \phi] &= 0 \quad \text{on } \Sigma, \\
\partial_{\nu} \phi &= 0 \quad \text{on } \partial \Omega,
\end{align*}
\]

by Theorem \[ 5.5 \]

Since \( \langle f | \nabla \phi \rangle = \langle u | \nabla \phi \rangle = 0 \) we obtain by integration by parts

\[
\langle \pi | g \rangle_{\Omega} = \left( \frac{\pi}{\rho} | \Delta \phi \rangle_{\Omega} = - \int_{\Omega} \left[ \frac{\pi}{\rho} \partial_{\nu} \phi \right] d\Sigma - \left( \frac{\mu}{\rho} \Delta u | \nabla \phi \rangle_{\Omega} \right. \\
\left. = - \int_{\Omega} \left[ \frac{\mu}{\rho} \nabla u : \nabla^2 \phi \right] dx + \int_{\Omega} \left\{ \frac{\mu}{\rho} \nabla u | \nabla \phi \rangle_{\Omega} - \left[ \frac{\mu}{\rho} \nabla \phi \right] d\Sigma. \right.
\]

Since \( \nabla u \in H^{1/2}(J; L_p(\Omega)^{n \times n}) \) and \([\pi], \partial_{\nu} u \in H^{1/2-1/2p}(J; L_p(\Sigma))\), applying \( \partial_{\alpha}^\alpha \) to this identity, we obtain the estimate

\[
|\partial_{\alpha}^\alpha \pi|_{H^{1/2}(J; \Omega)} \leq C \left\{ |\partial_{\alpha}^\alpha \nabla u|_{L_p(\Omega)} + |\partial_{\alpha}^\alpha |_{L_p(\Sigma)} + |\partial_{\alpha}^\alpha \partial_{\nu} u|_{L_p(\Sigma)} \right\},
\]

for each \( \alpha \in (0, 1/2 - 1/2p), \) hence \( \pi \in H^{1/2}(J; L_p(\Omega))\).

It is convenient to reduce the problem to the case

\[
u_0 = h_0 = f = f_d = u_{\Sigma} \cdot \nu_\Sigma = u_b \cdot \nu_{\partial \Omega} = 0.
\]

This can be achieved as follows. Suppose \((u, \pi, h)\) is a solution of \[ 5.1 \]. We introduce a further dummy variable \( q := [\pi] \); note that \( q \in Z_q := Y_{a^1} \). We decompose \( u = u_* + u_1, \pi = \pi_* + \pi_1, q = q_* + q_1, h = h_* + h_1 \) where

\[
h_*(t) = [2e^{-(t-\Delta \Sigma)^{1/2}} - e^{-2(t-\Delta \Sigma)^{1/2}}] h_0 + \]

\[
\left[ e^{-(t-\Delta \Sigma)^{1/2}} - e^{-2(t-\Delta \Sigma)^{1/2}} \right] (I - \Delta \Sigma)^{-1} \{ (u_0 | \nu_\Sigma) - (\nabla_\Sigma h) + g_0(0) \}, \quad t \geq 0.
\]

The function \( h_* \) belongs to \( Z_h \) and satisfies \( h_*(0) = h_0 \) and \( \partial_{\nu} h_*(0) = (u_0 | \nu_\Sigma) - (\nabla_\Sigma h) + g_0(0) \). Then \( h_* \) has initial value zero, and also \( \partial_{\nu} h_*(0) = 0 \). We set \( q_*(t) = e^{\Delta \Sigma^2} q_0 \) where

\[
q_0 := ([\mu(\nabla u_0 + [\nabla u_0]^T)] \nu_\Sigma | \nu_\Sigma) + \sigma \Delta \Sigma h_0 + (g(0) | \nu_\Sigma)
\]

is determined by the data, and we define \( \pi_* \) as the solution of

\[
\begin{align*}
\Delta \pi_* &= 0 \quad \text{in } \Omega \setminus \Sigma, \\
\partial_{\nu} \pi_* &= 0 \quad \text{on } \partial \Omega, \\
[\partial_{\nu} \pi_*] &= 0, \quad [\pi_*] = q_* \quad \text{on } \Sigma.
\end{align*}
\]

Note that \( q_* \in Y_{a^1} \), and \( \pi_* \in Z_{\pi} \), by Theorem \[ 5.5 \]. The function \( u_* \in Z_u \) is defined as the solution of the parabolic problem

\[
\begin{align*}
\rho \partial_t u_* - \mu \Delta u_* &= -\nabla \pi_* + \rho f \quad \text{in } \Omega \setminus \Sigma, \\
u_* &= u_b \quad \text{on } \partial \Omega, \\
[-\mu(\nabla u_* + [\nabla u_*]^T)] \nu_\Sigma &= g - q_* \nu_\Sigma + \sigma (\Delta \Sigma h_*) \nu_\Sigma \quad \text{on } \Sigma, \\
u(0) &= u_*, \quad \text{in } \Omega \setminus \Sigma,
\end{align*}
\]

Proof. Let \( g \in L_p'(\Omega) \) be given and solve the problem

\[
\begin{align*}
\Delta \phi &= \rho g \quad \text{in } \Omega \setminus \Sigma, \\
[\phi] &= 0 \quad \text{on } \Sigma, \\
[\rho^{-1} \partial_{\nu} \phi] &= 0 \quad \text{on } \Sigma, \\
\partial_{\nu} \phi &= 0 \quad \text{on } \partial \Omega,
\end{align*}
\]
which is uniquely solvable since the appropriate Lopatinskii-Shapiro conditions are satisfied; see [12]. Thus we may assume w.l.o.g. \( u_0 = h_0 = f = 0 \) and that the time traces of \( f_d, g \) and \( g_h \) are zero at time zero. Finally, to remove \( f_d \), we solve the transmission problem

\[
\begin{align*}
\Delta \psi &= f_d \quad \text{in} \quad \Omega \setminus \Sigma, \\
[\rho \psi] &= 0 \quad \text{on} \quad \Sigma, \\
[\partial_\nu \psi] &= 0 \quad \text{on} \quad \Sigma, \\
\partial_{\nu\Omega \Sigma} \psi &= 0 \quad \text{on} \quad \partial \Omega,
\end{align*}
\]

according to Theorem 8.6 where \( \tilde{f}_d := f_d - \text{div} \ u_s \). Since \( \partial \Omega \in C^3 \), the solution satisfies \( \nabla \psi \in Z_w \). Then setting \( u_2 = u_1 - \nabla \psi \) and \( \sigma_2 = \pi_1 + \rho \partial_\nu \psi - \mu \Delta \psi, \ h_2 = h_1 \), we see that we may assume also \( f_d = u_{\Sigma} \cdot \nu_{\Sigma} = u_h \cdot \nu_{\Omega \Sigma} = 0 \), the only non-vanishing data which remain are \( g, g_h, u_{\Sigma}, u_h \); note that the time traces at \( t = 0 \) of these functions are zero, and \( u_h \cdot \nu_{\Omega \Sigma} = 0 \) on \( \partial \Omega \) and \( u_{\Sigma} \cdot \nu_{\Sigma} = 0 \) on \( \Sigma \).

3.1. Flat Interface. In this subsection we consider the linearized problem for a flat interface.

\[
\begin{align*}
\rho \partial_t u - \mu \Delta u + \nabla \pi &= \rho f \quad \text{in} \quad \mathbb{R}^n, \\
\text{div} \ u &= f_d \quad \text{in} \quad \mathbb{R}^n, \\
- [\mu \partial_y v] - [\mu \nabla_x w] &= g_v \quad \text{on} \quad \mathbb{R}^{n-1}, \\
-2[\mu \partial_y w] + [\pi] - \sigma \Delta h &= g_w \quad \text{on} \quad \mathbb{R}^{n-1}, \\
[u] &= u_{\Sigma} \quad \text{on} \quad \mathbb{R}^{n-1}, \\
\partial_t h - w + (b \nabla h) &= g_h \quad \text{on} \quad \mathbb{R}^{n-1}, \\
u(0) &= u_0, \ h(0) = h_0 \quad \text{in} \quad \mathbb{R}^n, \ \text{on} \ \mathbb{R}^{n-1}. \tag{3.4}
\end{align*}
\]

Here we have identified \( \mathbb{R}^{n-1} = \mathbb{R}^{n-1} \times \{0\} \) and \( \mathbb{R}^n = \mathbb{R}^n \times \mathbb{R}^{n-1} \). It is convenient to split \( u = (v, w) \), \( f = (f_v, f_w) \), \( g = (g_v, g_w) \) into tangential and normal components.

The following result, which is implied by [20, Theorem 3.1], states that problem (3.4) admits maximal regularity, in particular defines an isomorphism between the solution space \( Z := Z_u \times Z_\pi \times Z_q \times Z_h \) and the product-space of data \((u_0, h_0, u_\Sigma, f, f_d, g, g_h, b)\) which we denote for short by \( Y \).

**Proposition 3.3.** Let \( p > n + 2 \) be fixed, and assume that \( \rho_j, \mu_j, \sigma \) are positive constants for \( j = 1, 2 \), and let \( J = [0, a] \). Suppose

\[
b_0 \in \mathbb{R}^{n-1}, \quad b_1 \in W^{1-1/p}_p(J; L_p(\mathbb{R}^{n-1}))^{n-1} \cap L_p(J; W^{2-1/p}_p(\mathbb{R}^{n-1}))^{n-1},
\]

and set \( b = b_0 + b_1 \). Then the Stokes problem with flat boundary (3.4) admits a unique solution

\[
(u, \pi, h) \quad \text{with regularity}
\]

\[
\begin{align*}
u &\in H^1_p(J; L_p(\mathbb{R}^n)^n) \cap L_p(J; H^2_p(\mathbb{R}^n)^n), \quad \pi \in L_p(J; H^1_p(\mathbb{R}^n)), \\
[\pi] &\in W^{1/2-1/p}_p(J; L_p(\mathbb{R}^{n-1})) \cap L_p(J; W^{1/p}_p(\mathbb{R}^{n-1})), \\
h &\in W^{2-1/p}_p(J; L_p(\mathbb{R}^{n-1})) \cap H^1_p(J; W^{2-1/p}_p(\mathbb{R}^{n-1})) \cap L_p(J; W^{2-1/p}_p(\mathbb{R}^{n-1})),
\end{align*}
\]

if and only if the data \((f, f_d, g, g_h, u_0, h_0, u_\Sigma)\) satisfy the following regularity and compatibility conditions:
We may assume by the reduction explained above that

\[ u \in W^{1,2/p}(J; L^p(\mathbb{R}^n)) \cap L^p(J; W^{2-1/p}(\mathbb{R}^n)), \]

\[ f_d \in L^p(J; H^1(\mathbb{R}^n)), \quad (f_d, u^\Sigma) \in H^1(J; \tilde{\mathcal{H}}^{-1}(\mathbb{R}^n)), \]

\[ g = (g_v, g_w) \in W^{1,2/p}(J; L^p(\mathbb{R}^n)) \cap L^p(J; W^{2-1/p}(\mathbb{R}^n)), \]

\[ g_h \in W^{1,2/p}(J; L^p(\mathbb{R}^n)) \cap L^p(J; W^{2-1/p}(\mathbb{R}^n)), \]

\[ u_0 \in W^{2-2/p}(\mathbb{R}^n), \quad h_0 \in W^{3-2/p}(\mathbb{R}^n), \]

\[ \text{div } u_0 = f_d|_{t=0} \text{ in } \mathbb{R}^n \text{ and } \|u_0\| = u_\Sigma|_{t=0} \text{ on } \mathbb{R}^n, \]

\[ \text{div } u_0 = g_d|_{t=0} \text{ on } \mathbb{R}^n. \]

The solution map \([u_0, h_0, u^\Sigma, f, f_d, g, g_h, b] \mapsto (u, \pi, [\pi], h)\) is continuous between the corresponding spaces.

### 3.2. Bent Interfaces

Next we consider the case of a bent interface. By this we mean a situation where the interface \(\Sigma\) is given as a graph of a function \(\phi : \mathbb{R}^{n-1} \to \mathbb{R}\) of class \(BC^3\); thus \(\Sigma = \{(x, \phi(x)) : x \in \mathbb{R}^{n-1}\}\). The normal \(\nu^\Sigma\) is then given by

\[ \nu^\Sigma(x) = \beta(x) \left[ -\nabla x \phi(x) \overline{1} \right], \quad \beta(x) = 1 / \sqrt{1 + |\nabla x \phi(x)|^2}, \]

and the Laplace-Beltrami operator for such a surface with \(\tilde{h}(t, x) = h(t, (x, \phi(x)))\) reads as

\[ \Delta^\Sigma h = \Delta \tilde{h} - \beta^2 (\nabla^2 \tilde{h} \nabla \phi | \nabla \phi) - \beta^2 (\Delta \phi - \beta^2 (\nabla^2 \phi \nabla \phi | \nabla \phi)) (\nabla \phi | \nabla \tilde{h}). \]

We may assume by the reduction explained above that \(u_0 = h_0 = f = f_d = u^\Sigma \cdot \nu^\Sigma = 0\). Set

\[ \tilde{u}(t, x, y) = u(t, x, y + \phi(x)), \quad \tilde{\pi}(t, x, y) = \pi(t, x, y + \phi(x)), \]

for \(t \in J = [0, a], x \in \mathbb{R}^{n-1}, y \neq 0\), and observe

\[ \nabla u = \nabla \tilde{u} - \nabla \phi \otimes \partial_y \tilde{u}. \]

Then we obtain for the new variables \((\tilde{u}, \tilde{\pi}, \tilde{h})\) the following problem. For convenience we drop the bars, and split \(u = (v, w)\) and \(g = (g_v, g_w)\) as before.

\[ \rho \partial_t v - \mu \Delta v + \nabla \pi = \mu B_1 (u, \pi), \quad \text{in } \mathbb{R}^n, \]

\[ \text{div } u = B_2 u, \quad \text{in } \mathbb{R}^n, \]

\[ -[\mu \partial_y v] - [\mu \nabla_x w] = g_v + B_3(u, [\pi], h), \quad \text{on } \mathbb{R}^{n-1}, \]

\[ -2[\mu \partial_y w] + [\pi] - \sigma \Delta_x h = [g_w / \beta] + B_4(u, h), \quad \text{on } \mathbb{R}^{n-1}, \quad (3.5) \]

\[ [u] = u^\Sigma, \quad \text{on } \mathbb{R}^{n-1}, \]

\[ \partial_t h - w + (b | \nabla h) = g_h + B_5 u + B_6 h, \quad \text{on } \mathbb{R}^{n-1}, \]

\[ u(0) = 0, \quad h(0) = 0, \quad \text{in } \mathbb{R}^n, \text{ on } \mathbb{R}^{n-1}. \]
Here we have set

\[ B_1(u, \pi) = |\nabla \phi|^2 \partial^2_{\nu u} u - 2(\nabla \phi | \nabla_x \partial_y u) + (\nabla \phi) \partial_y \pi - (\Delta \phi) \partial_y u \]
\[ B_2 u = (\nabla \phi) \partial_y u, \]
\[ B_3(u, [\pi], h) = -[\mu(\nabla_x v + [\nabla_x v]^T)] \nabla \phi - [\mu \partial_y v] | \nabla \phi |^2 \]
\[ + \{-[\mu(\partial_y v | \nabla \phi)] + [\pi] - [\mu \partial_y w] - \sigma \Delta \phi \} \nabla \phi \]
\[ B_4(u, h) = -([\mu(\partial_y v + \nabla_x w)] | \nabla \phi | - [\mu \partial_y w] | \nabla \phi |^2 + \sigma(\Delta \phi h - \Delta h) \]
\[ B_5 u = (\beta - 1) w - \beta(\nabla \phi | v) = -\frac{\beta^2 | \nabla \phi |^2}{1 + \beta} - w - \beta(\nabla \phi | v) \]
\[ B_6 h = \beta^2 [(b | \nabla \phi | - (b \epsilon_n) | \nabla \phi |^2)(\nabla \phi | \nabla h)]. \]

Now suppose \((u, \pi, h)\) belongs to the maximal regularity class. We estimate the perturbations \(B_j\) as follows:

\[ |B_1(u, \pi)|_{L^p} \leq ||\nabla \phi||_{L^\infty} [(2 + |\nabla \phi|_{L^\infty}) |\nabla^2 u|_{L^p} + |\nabla \pi|_{L^p}] \]
\[ + |\Delta \phi|_{L^\infty} |\nabla u|_{L^p}, \]
\[ |B_2 u|_{L^p(H^1_p)} \leq |\nabla \phi|_{L^\infty} |\nabla^2 u|_{L^p} + (|\nabla^2 \phi|_{L^\infty} + |\nabla \phi|_{L^\infty}) |\nabla u|_{L^p}, \]
\[ |\partial_y B_2 u|_{L^p(H^{-1}_p)} \leq |\nabla \phi \partial_y u|_{L^p} \leq |\nabla \phi|_{L^\infty} |\partial_y u|_{L^p}, \]
\[ |B_3(u, [\pi], h)|_{L^p(H^1_p)} \leq C |\nabla \phi|_{L^\infty} (1 + |\nabla \phi|_{L^\infty}) |\nabla u|_{W^1_p(L^p)} + |\nabla^2 h|_{W^1_p(L^p)} \]
\[ + C |\nabla \phi|_{L^\infty} ([\nabla \pi]_{W^2_p(L^p)} + |\nabla^2 \phi|_{L^\infty} |\nabla h|_{W^2_p(L^p)}), \]
\[ |B_4(u, h)|_{W^1_p(L^p)} \leq C |\nabla \phi|_{L^\infty} (1 + |\nabla \phi|_{L^\infty}) |\nabla u|_{W^1_p(L^p)} + |\nabla^2 h|_{W^1_p(L^p)} \]
\[ + C |\nabla^2 \phi|_{L^\infty} |\nabla \phi|_{L^\infty} |\nabla h|_{W^1_p(L^p)} \]
\[ |B_5 u|_{W^{1-1/2p}_p(L^p)} \leq 2 |\nabla \phi|_{L^\infty} |\nabla u|_{W^{1-1/2p}_p(L^p)}. \]

Here \(C\) denotes a constant only depending on the parameters \(\mu\) and \(\sigma\), and we have set \(s = 1/2 - 1/2p\). For the estimations in \(L^p(J; W^{1-1/p}_p(\mathbb{R}^{n-1}))\) we observe that

\[ |\psi|_{W^s_p} = |\psi|_{L^p} + [\psi]_{x, p}, \quad [\psi]^p_{x, p} = \int_{|h| \leq 1} \int_{\mathbb{R}^{n-1}} |\psi(x + h) - \psi(x)|^p \frac{dx \, dh}{|h|^{n-1+sp}}, \]

defines a norm on \(W^s_p(\mathbb{R}^{n-1})\). This implies

\[ |a \psi|_{W^s_p} \leq |a|_{L^\infty} |\psi|_{W^s_p} + c_{s, p} |\nabla u|_{L^\infty} |\psi|_{L^p}, \]
for \( a \in W^1_2(\mathbb{R}^{n-1}) \), with some constant \( c_{s,p} \) which only depends on \( s \in (0, 1) \), \( p \in (1, \infty) \) and on \( n \). With this observation we have

\[
|B_3(u, [\pi], h)|_{L_p(W^{2s}_p)} \leq C(1 + |\nabla \phi|_{L\infty})^i |\nabla \phi|_{L\infty}:
\]

\[
\cdot \left( |\nabla u|_{L_p(W^{2s}_p)} + |[\pi]|_{L_p(W^{2s}_p)} + |\nabla h|_{L_p(H^{1s}_p)} \right) + |\nabla \phi|_{L\infty} \left( |\nabla u|_{L_p} + |[\pi]|_{L_p} + |\nabla h|_{L_p(H^{1s}_p)} \right)
\]

\[
+ C(|\nabla^2 \phi|_{L\infty} + |\nabla^2 \phi|_{L\infty}) |\nabla \phi|_{L\infty} |\nabla h|_{L_p},
\]

\[
|B_4(u, h)|_{L_p(W^{2s}_p)} \leq C(1 + |\nabla \phi|_{L\infty})^i |\nabla \phi|_{L\infty} |\nabla u|_{L_p(W^{2s}_p)}
\]

\[
+ |\nabla \phi|_{L\infty} \left( |\nabla^2 h|_{L_p(W^{2s}_p)} + |\nabla^2 \phi|_{L\infty} \left( |\nabla u|_{L_p} + |\nabla h|_{L_p(H^{1s}_p)} \right) \right)
\]

\[
+ C(|\nabla^2 \phi|_{L\infty} + |\nabla^2 \phi|_{L\infty}) |\nabla \phi|_{L\infty} |\nabla h|_{L_p},
\]

\[
|B_5u|_{L_p(W^{2s}_p)} \leq C|\nabla \phi|_{L\infty} \left\{ |u|_{L_p(W^{1s}_p)} + |\nabla \phi|_{L\infty} \right\} |u|_{L_p(H^{1s}_p)}
\]

\[
+ C(|\nabla^3 \phi|_{L\infty} + |\nabla^2 \phi|_{L\infty}) |u|_{L_p},
\]

Here \( C \) denotes a constant only depending on \( \mu, \sigma, p, \) and \( 2s = 1 - 1/p \). To estimate \( B_0h \) we note that \( Y^a_u \) is a Banach algebra since \( p > n + 2 \). This yields

\[
|B_0h|_{Y^a_u} \leq C|\nabla \phi|_{L\infty} \left\{ |b_0| + |b_1|_{Y^a_u} \right\} |\nabla h|_{Y^a_u} \leq C|\nabla \phi|_{L\infty} \left\{ |b_0| + |b_1|_{Y^a_u} \right\} |h|_{Zp}.
\]

To solve the problem (3.3), let \( z = (u, \pi, [\pi], h) \in \partial Z \), where \( \partial Z \) means the solution space with zero time trace at \( t = 0 \). \( F := (0, 0, g_u, g_v, \beta, u_\Sigma, g_h) \in \partial Y \), the space of data with zero time trace, and let \( B : \partial Z \to \partial Y \) be defined by

\[
Bz = (B_1(u, \pi), B_2u, B_3(u, [\pi], h), B_4(u, h), 0, B_5u + B_0h).
\]

Denoting the isomorphism from \( \partial Z \) to \( \partial Y \) by the left hand side of (3.3) by \( L \) we may rewrite problem (3.3) in abstract form as

\[
Lz = Bz + F. \tag{3.6}
\]

The above estimates for the components of \( B \) imply

\[
|Bz|_Y \leq C|\nabla \phi|_{L\infty} |z|_Z + M \left\{ |u|_{L_p(H^{1s}_p)} + |[\pi]|_{L_p} + |\nabla h|_{W^{2s}_p(\Sigma)\cap L_p(H^{1s}_p)} \right\},
\]

with some constants \( C > 0 \) depending only on the parameters and \( M > 0 \), which depends also on \( |\nabla \phi|_{BLUC^2} \). Let \( \eta > 0 \) be given and suppose \( |\nabla \phi|_{L\infty} < \eta \). By means of an interpolation argument we find a constant \( \eta > 0 \), depending only on \( p \) such that there is a constant \( M(\eta) > 0 \) such that

\[
|Bz|_Y \leq C|2\eta + a'M(\eta)||z|_Z, \quad z \in \partial Z.
\]

Choosing first \( \eta > 0 \) and then \( a > 0 \) small enough, we can solve (3.6) by a Neumann series argument for \( J = [0, a] \).

Since problem (3.3) is time-invariant, we may repeat these arguments finitely many times, including the reduction procedure, to solve (3.3) for \( J = [0, a] \), where now \( a > 0 \) is arbitrary.

3.3. **General Bounded Geometries.** Here we use the method of localization. By assumption, \( \partial Z \) is of class \( C^3 \) and \( \Sigma \) will even be real analytic, so in particular of class \( C^4 \). Therefore we may cover \( \Sigma \) by \( N \) balls \( B_{r/2}(x_j) \) with radius \( r > 0 \) and centers \( x_j \in \Sigma \) such that \( \Sigma \cap B_{r}(x_j) \) can be parameterized over the tangent space \( T_{x_j}\Sigma \) by a function \( \theta_j \in C^4 \) such that \( |\nabla \theta_j|_{L\infty} \leq \eta \), with \( \eta > 0 \) defined as in the previous subsection. We extend these functions \( \theta_j \) to all of \( T_{x_j}\Sigma \), retaining the bound on \( \nabla \theta_j \). This way we have created \( N \) bent half-spaces \( \Sigma_j \) to which the
result proved in the previous subsection applies. We also suppose that $B_r(x_j) \subset \Omega$ for each $j$. Set $U := \Omega \cup_{j=1}^N B_{r/2}(x_j)$ and $U_j = B_r(x_j)$, $j = 1, \ldots, N$. The open set $U$ consists of one component $U_0$ characterized by $\partial \Omega \subset U_0$ and an open set say $U_{N+1}$, which is interior to $\Sigma$, i.e. $U_j \subset \Omega_1$. Fix a partition of unity $\{\varphi_j\}_{j=0}^{N+1}$ subject to the covering $\{U_j\}_{j=0}^{N+1}$ of $\Omega$, i.e. $\varphi_j \in \mathcal{D}(\mathbb{R}^n)$, $0 \leq \varphi_j \leq 1$, and $\sum_{j=0}^{N+1} \varphi_j \equiv 1$. Note that $\varphi_0 = 1$ in a neighborhood of $\partial \Omega$. Let $\hat{\varphi}_j$ denote cut-off functions with support in $U_j$ such that $\hat{\varphi}_j = 1$ on the support of $\varphi_j$, and set $b_j = b \hat{\varphi}_j$.

Let $z := (u, \pi, q, h)$ with $q = [\pi]$ be a solution of (3.6) where we assume w.l.o.g. $u_0 = h_0 = h_1 = f = f_d = 0$, and $(u_b|_{\partial \Omega}) = [(u_{\Sigma}|_{\Sigma})] = 0$. We then set $u_j = \varphi_j u$, $\pi_j = \varphi_j \pi$, $q_j = \varphi_j q$, $h_j = \varphi_j h$, as well as $u_{b,j} = \varphi_j u_0$, $u_{\Sigma,j} = \varphi_j u_{\Sigma}$, $g_j = \varphi_j g$, and $g_{b,j} = \varphi_j g_h$. Then for $j = 1, \ldots, N$, the quadruples $z_j := (u_j, \pi_j, q_j, h_j)$ satisfy the problems

$$\begin{align*}
\rho \partial_t u_j - \mu \Delta u_j + \nabla \pi_j &= F_j(u, \pi) \quad \text{in } \mathbb{R}^n \setminus \Sigma_j, \\
\text{div } u_j &= (\nabla \varphi_j | u) \quad \text{in } \mathbb{R}^n \setminus \Sigma_j, \\
[-\mu((\nabla u_j) + [\nabla u_j]^T) + q_j] \nu_{\Sigma_j} - \sigma(\Delta \nu_{\Sigma_j} h_j) \nu_{\Sigma_j} &= g_j + G_j(u) \quad \text{on } \Sigma_j, \\
[u_j] = u_{\Sigma,j}, \quad g_j &= [\pi_j] \quad \text{on } \Sigma_j, \\
\partial_t h_j - (u_j | \nu_{\Sigma_j}) + (b_j | \nabla \nabla h_j) &= g_{b,j} + G_{h,j}(h) \quad \text{on } \Sigma_j, \\
u_j(0) &= 0 \quad \text{in } \mathbb{R}^n \setminus \Sigma_j, \quad h_j(0) = 0 \quad \text{on } \Sigma_j.
\end{align*}$$

(3.7)

Here we used the abbreviations

$$\begin{align*}
F_j(u, \pi) &= [\nabla \varphi_j | \pi - \mu \Delta \varphi_j] u, \\
G_j(u) &= [-\mu((\nabla \varphi_j \otimes u + u \otimes \nabla \varphi_j) | \nu_{\Sigma_j}) - \sigma(\Delta \varphi_j h) \nu_{\Sigma_j},
\end{align*}$$

and

$$G_{h,j}(h) = (b_j | \nabla \nabla \varphi_j) h.$$ 

For $j = 0$ we have the standard one-phase Stokes problem with parameters $\rho_2, \mu_2$ on $\Omega$ with Dirichlet boundary conditions on $\partial \Omega$, i.e.

$$\begin{align*}
\rho_2 \partial_t u_0 - \mu_2 \Delta u_0 + \nabla \pi_0 &= F_0(u, \pi) \quad \text{in } \Omega, \\
\text{div } u_0 &= (\nabla \varphi_0 | u), \quad \text{in } \Omega, \\
\quad u_0 &= u_{\lambda,0} \quad \text{on } \partial \Omega, \\
\quad u_0(0) &= 0 \quad \text{in } \Omega.
\end{align*}$$

For $j = N + 1$ we obtain the one-phase Stokes problem on $\mathbb{R}^n$ with parameters $\rho_1, \mu_1$, i.e.

$$\begin{align*}
\rho_1 \partial_t u_{N+1} - \mu_1 \Delta u_{N+1} + \nabla \pi_{N+1} &= F_{N+1}(u, \pi) \quad \text{in } \mathbb{R}^n, \\
\text{div } u_{N+1} &= (\nabla \varphi_{N+1} | u) \quad \text{in } \mathbb{R}^n, \\
\quad u_{N+1}(0) &= 0 \quad \text{in } \mathbb{R}^n.
\end{align*}$$

Concentrating on $j = 1, \ldots, N$, we first note that $[\Delta, \varphi_j]$ are differential operators of order 1, hence if $u \in 0 Z_u$ then

$$[\Delta, \varphi_j] u \in 0 H^{1/2}_p(J; L_p(\mathbb{R}^n)^n) \cap L_p(J; H^1_p(\mathbb{R}^n\setminus \Sigma_j)^n).$$

Since $f = f_d = 0$ the pressure $\pi$ belongs to

$$\pi \in 0 H^0_p(J; L_p(\mathbb{R}^n)) \cap L_p(J; H^1_p(\mathbb{R}^n\setminus \Sigma_j)),$$

by Corollary 3.2 hence we have

$$F_j(u, \pi) \in 0 H^0_p(J; L_p(\mathbb{R}^n)^n) \cap L_p(J; H^1_p(\mathbb{R}^n\setminus \Sigma_j))^n,$$
for some fixed $0 < \alpha < \frac{1}{2} - \frac{1}{2p}$. Similarly we have
\[
\nabla \varphi_j (u |_{\Sigma_j}) + (\nabla \varphi_j |_{\Sigma_j}) u \in \omega W^{\frac{1}{2}, 2p}(J; L_p(\Sigma_j)^n) \cap L_p(J; W^{2-1/p}(\Sigma_j)^n),
\]
and since $[\Delta_{\Sigma_j}, \varphi_j]$ is of order 1 as well, we obtain
\[
[\Delta_{\Sigma}, \varphi_j] h \in \omega H^1(J; W^{1-1/p}(\Sigma_j)) \cap L_p(J; W^{2-1/p}(\Sigma_j)).
\]
This shows that we have
\[
G_j(u) \in \omega W^{1-1/p}(J; L_p(\Sigma_j)^n) \cap L_p(J; W^{2-1/p}(\Sigma_j)^n).
\]

The terms $G_{h_j}(h)$ do not have more regularity, however, the Banach algebra property yields the estimate
\[
|G_{h_j}(h)|_{Y_0^1} \leq C |b|_{Y_0^1} \|h|_{Y_0^1} \leq C |b|_{Y_0^1}^{a^2} \|h|_{Z_a},
\]
with an appropriate exponent $\gamma > 0$. Next we decompose
\[
F_j(u, \pi) = \tilde{F}_j(u, \pi) + \nabla \psi_j,
\]
such that $\text{div} \tilde{F}_j(u, \pi) = 0$ in $\mathbb{R}^n \setminus \Sigma_j$ and $([\tilde{F}_j(u, \pi)] |_{\Sigma_j}) = 0$ on $\Sigma_j$. Thus $\tilde{F}_j(u, \pi)$ is the Helmholtz projection of $F_j(u, \pi)$ in $\mathbb{R}^n$. Then
\[
\tilde{F}_j(u, \pi) \in \omega H^1_p(J; L_p(\mathbb{R}^n)^n) \cap L_p(J; H^{2\alpha}_p(\mathbb{R}^n)^n).
\]

Also, we decompose $u_j = \tilde{u}_j + \nabla \phi_j$, where $\phi_j$ solves the transmission problem
\[
\nabla \phi_j \in \omega H^1(\Sigma; H^1(\mathbb{R}^n \setminus \Sigma_j)^n) \cap L_p(\Sigma; H^{2\alpha}(\mathbb{R}^n \setminus \Sigma_j)^n),
\]
and its normal part vanishes, by construction. Further we have
\[
[\mu \nabla^2 \phi_j] \in \omega W^{1-1/2p}(J; L_p(\Sigma_j)^n \times \Sigma_j) \cap L_p(J; W^{2-1/p}(\Sigma_j)^n \times \Sigma_j).
\]
Then we set
\[
\tilde{\pi}_j = \pi_j - \psi_j + \rho \partial_\nu \phi_j - \mu \Delta \phi_j,
\]
and we observe that on $\Sigma_j$
\[
\tilde{q}_j := [\tilde{\pi}_j] = [\pi_j] - [\mu \Delta \phi_j] = [\pi_j] - [\mu (\nabla \varphi_j | u)],
\]
which by construction $\psi_j$ and $\rho \phi_j$ have no jump across $\Sigma_j$. Now the quadruple $\tilde{z}_j := (\tilde{u}_j, \tilde{\pi}_j, \tilde{q}_j, h_j)$ satisfies the problem
\[
\rho \partial_\nu \tilde{u}_j - \mu \Delta \tilde{u}_j + \nabla \tilde{\pi}_j = \tilde{F}_j(u, \pi) \quad \text{in } \mathbb{R}^n \setminus \Sigma_j,
\]
\[
\text{div } \tilde{u}_j = (\nabla \varphi_j | u) \quad \text{in } \mathbb{R}^n \setminus \Sigma_j,
\]
\[
[-\mu((\nabla \tilde{u}_j) + (\nabla \tilde{u}_j)^T) + \tilde{q}_j] \nu_{\Sigma_j} - \sigma (\Delta_{\Sigma_j} h_j) \nu_{\Sigma_j} = g_j + \tilde{G}_j(u) \quad \text{on } \Sigma_j,
\]
\[
[\tilde{u}_j] = u_{\Sigma_j} - [\nabla \phi_j], \quad \tilde{q}_j = [\tilde{\pi}_j] \quad \text{on } \Sigma_j,
\]
\[
\partial_\nu h_j - (\tilde{u}_j |_{\Sigma_j}) + (b_j | \nabla \phi_j) = g_{h_j} + \tilde{G}_{h_j}(h) \quad \text{on } \Sigma_j,
\]
\[
\tilde{u}_j(0) = 0, \quad \text{in } \mathbb{R}^n \setminus \Sigma_j, \quad h_j(0) = 0 \quad \text{on } \Sigma_j.
\]
Here $\tilde{G}_j$ and $\tilde{G}_{hj}$ are given by
\[
\tilde{G}_j(u) = G_j(u) + 2[\mu \nabla^2 \phi_j]|_{\Sigma_j} - [\mu(\nabla \varphi_j|u)]|_{\Sigma_j},
\]
and
\[
\tilde{G}_{hj}(h) = G_{hj}(h) + \partial_{\Sigma_j} \phi_j.
\]
For the remaining charts with index $j = 0, N + 1$, i.e. the one-phase problems, the procedure is similar. In case $j = 0$ we use the regularity
\[
\nabla \phi_0 \in H^1_p(J; H^1_p(\Omega)^n) \cap H^{1/2}_p(J; H^2_p(\Omega)^n)
\]
instead of (3.8).
We write (3.9) abstractly as
\[
L_j \tilde{z}_j = H_j + B_j z,
\]
and by Theorem 3.1 for bent interfaces we obtain an estimate of the form
\[
|\tilde{z}_j|_E \leq C_0(|H_j|_F + |B_j z|_F),
\]
with some constant $C_0$ independent of $j$. Here $E$ means the space of solutions and $F$ the space of data. Since all components of $B_j z$ (except for $G_{hj}(h)$) have some extra regularity, there is an exponent $\gamma > 0$ and a constant $C_1$ independent of $j$ such that
\[
|B_j z|_F \leq a^{\gamma} C_1 |z|_E.
\]
In addition, by Corollary 3.2 we obtain
\[
|\tilde{\gamma}_j|_{H^1_p(J; L_p(\Omega))} \leq C_2(|H_j|_F + |B_j z|_F) \leq C_2(|H_j|_F + a^{\gamma} C_1 C_2 |z|_E).
\]
This in turn implies
\[
|\partial_\nu \phi_j|_{H^1_p(J; L_p(\Omega))} \leq C_2 |H_j|_F + a^{\gamma} C_3 |z|_E,
\]
and then also
\[
|z_j|_E \leq C_4 |H_j|_F + a^{\gamma} C_5 |z|_E.
\]
Summing over all $j$ yields $z = \sum_j z_j$, hence
\[
|z|_E \leq C_6 |H|_F + a^{\gamma} C_7 |z|_E.
\]
Therefore, choosing the length $a$ of the time interval small enough, we obtain the a priori estimate
\[
|z|_E \leq C_8 |H|_F. \tag{3.10}
\]
Since the equations under consideration are time invariant, repeating this argument finitely many times we may conclude that the operator $L : aE \to aF$ which maps solutions to their data is injective and has closed range, i.e. $L$ is a semi-Fredholm operator.

It remains to prove surjectivity of $L$. For this we employ the continuation method for semi-Fredholm operators. The estimates are uniform in the densities $\rho_j$ and the viscosities $\mu_j$, as long as these parameters are bounded and bounded away from zero. Hence $L = L(\rho_1, \rho_2, \mu_1, \mu_2)$ is surjective, if $L(1, 1, 1, 1)$ has this property. Next we introduce an artificial continuation parameter $\tau \in [0, 1]$ by replacing the equation for the free boundary $h$ with
\[
\partial \tau h + (1/2) h - (1 - \tau)\{(u|_{\Sigma}) - (b|_{\Sigma} h)\} = g_h \quad \text{on } \Sigma.
\]
The arguments in [10] [20] show that the corresponding problem is well-posed for each $\tau \in [0, 1]$ in the case of a flat interface, with bounds independent of $\tau \in [0, 1]$. 
Therefore the same is true for bent interfaces and then by the above estimates also for a general geometry. Thus we only need to consider the case \( \rho_1 = \rho_2 = \mu_1 = \mu_2 = \tau = 1 \).

To prove surjectivity in this case, note that the equation for \( h \) is decoupled from those for \( u \) and \( \pi \), and it is uniquely solvable in the right regularity class because of maximal regularity for the Laplace-Beltrami operator. So we may set now \( h = 0 \). Next we solve the parabolic transmission problem to remove the jump of \( u \) across \( \Sigma \) and the inhomogeneity \( g \) for a general geometry. Thus we only need to consider the case \( \rho_1 = \rho_2 = \mu_1 = \mu_2 = \tau = 1 \), hence also for arbitrary \( \rho, \mu \) and \( \tau = 0 \) and the proof of Theorem 3.1 is complete.

We close this section with a remark on the situation, which occurs in the treatment of two-phase flows with variable surface tension. To be precise, Theorem 3.1 may be generalized to this situation by means of the following

**Corollary 3.4.** The statements of Theorem 3.1 remain valid, if the surface tension is not a constant but a function \( \sigma \in C^{0,1/2}(J,BC(\Sigma,\mathbb{R}^+)) \cap BC(J, C^{0,1}(\Sigma,\mathbb{R}^+)) \).

This generalization is possible, since the variability of the surface tension may be handeled by a freezing technique during the localization procedure in the previous paragraphs. The only difference compared to the case of a constant surface tension occurs in the jump condition of the normal stress in problem (3.7), which has to be modified to

\[
\left[ -\mu(\nabla u_j + [\nabla u_j]^T)\right] \nu_{\Sigma_j} + q_j \nu_{\Sigma_j} - \sigma_j(\Delta_{\Sigma_j} h_j) \nu_{\Sigma_j} = g_j + G_j(u) + G_{\sigma_j}(h) \quad \text{on} \ \Sigma_j
\]

with \( \sigma_j = \sigma(0, x_j) > 0 \) and

\[
G_{\sigma_j}(h) = \varphi_j(\sigma - \sigma_j)(\Delta_{\Sigma_j} h) \nu_{\Sigma_j}.
\]

Now, if the radius \( r > 0 \) of the balls \( U_j \), \( j = 1, \ldots, N \) is sufficiently small, we may use the Hölder/Lipschitz constants \( L_1, L_2 \) of \( \sigma \) to estimate

\[
|\varphi_j(\sigma - \sigma_j)|_{Y_2^1} \leq C_9(L_1 a^{\gamma_1} + L_2 r^{\gamma_2})
\]

with exponents \( \gamma_1, \gamma_2 > 0 \) and a constant \( C_9 > 0 \) depending only on the spatial dimension \( n \) and the geometry of \( \Sigma \). Since \( Y_2^1 \) is a Banach algebra, we obtain

\[
|G_{\sigma_j}(h)|_{Y_2^1} \leq C_{10}\varphi_j(\sigma - \sigma_j)|_{Y_2^1} \Delta_{\Sigma_j} h|_{Y_2^1} \leq C_{11}(a^{\gamma_1} + r^{\gamma_2})|z|_G
\]

and the a priori estimate (3.11) remains valid, if both \( a \) and \( r \) are chosen sufficiently small.

4. Local Well-Posedness

We now turn to problem (1.1), and show its local well-posedness for given initial data \( \Gamma_0 \in W_p^{2-2/p} \) and \( u_0 \in W_p^{2-2/p}(\Omega \setminus \Gamma_0)^n \), which are subject to the compatibility conditions

\[
\frac{\partial u_0}{\partial \nu} = 0 \text{ in } \Omega \setminus \Gamma_0, \quad u_0 = 0 \text{ on } \partial \Omega,
\]

\[
[P_{\Gamma_0} \mu(\nabla u_0 + [\nabla u_0]^T) \nu_{\Gamma_0}] = 0, \quad [u_0] = 0 \quad \text{on} \ \Gamma_0,
\]

where \( P_{\Gamma_0} = I - \nu_{\Gamma_0} \times \nu_{\Gamma_0} \). According to the considerations in Section 2 we approximate \( \Gamma_0 \) for any prescribed \( \eta > 0 \) by a real analytic hypersurface \( \Sigma \), in the sense that
We take \( d_H(N^2 \Sigma, N^2 \Gamma_0) < \eta \), and \( \Gamma_0 \) is parametrized over \( \Sigma \) by \( h_0 \in W_p^{3-2/p}(\Sigma) \). Employing the transformation from Section 2 to the fixed domain, it is sufficient to prove the local well-posedness of the quasilinear problem \((2.2)\). We keep the notation and denote by \( u \) and \( \pi \) the transformed velocity field and pressure, respectively.

We are interested in solutions of \((2.2)\) having maximal regularity, and hence, we determine a first approximation of the local solution of \((2.2)\) by using Theorem 3.1.

Since the time traces

\[
(\nabla | u_0 |) \in W_p^{1-2/p}(\Omega^t) \quad \text{and} \quad P_\Sigma(\mu | \nabla u_0 | [T]) \in W_p^{1-3/p}(\Sigma)
\]

will not be trivial due to the transformation, and since the compatibility conditions imposed in Theorem 3.1 have to be satisfied, we must be able to construct extensions to the right regularity classes to be used as the right-hand sides in \((3.1)\).

**Proposition 4.1.** Let \( p > 3 \), \( \partial \Omega \in C^3 \), and set \( J = [0, a] \). Let \( \Sigma \subset \Omega \) be a closed hypersurface of class \( C^3 \). Then \( f_0 \in \dot{H}_p^{-1}(\Omega) \cap W_p^{1-2/p}(\Omega^t) \) and \( g_0 \in W_p^{1-3/p}(\Sigma) \) admit extensions

\[
f \in H_p^1(J; \dot{H}_p^{-1}(\Omega^t)) \cap L_p(J; H_p^1(\Omega^t)) \quad \text{and} \quad g \in W_p^{1-2/p}(J; L_p(\Sigma^t)) \cap L_p(J; W_p^{1-1/p}(\Sigma))
\]

with \( f(0) = f_0 \) and \( g(0) = g_0 \).

**Proof.** We take \( \phi_0 \in H_p^2(\Omega) \cap W_p^{3-2/p}(\Omega^t) \) to be the unique solution of

\[
\begin{align*}
\Delta \phi_0 &= f_0 \quad \text{in} \quad \Omega^t, \\
\| \phi_0 \| &= 0 \quad \text{on} \quad \Sigma, \\
[\partial_\nu \phi_0] &= 0 \quad \text{on} \quad \Sigma, \\
\phi_0 &= 0 \quad \text{on} \quad \partial \Omega,
\end{align*}
\]

which exists due to \( f_0 \in \dot{H}_p^{-1}(\Omega^t) \cap W_p^{1-2/p}(\Omega^t) \); cf. Theorems 8.5 and 8.6. Setting \( v_0 := \nabla \phi_0 \in W_p^{2-2/p}(\Omega^t)^n \), the parabolic problem

\[
\begin{align*}
\partial_t v - \Delta v &= 0 & \text{in} \ J \times \Omega^t, \\
\{ v \} &= 0 & \text{on} \ J \times \Sigma, \\
[\partial_\nu v] &= e^{\Delta \Sigma t} [\partial_\nu v_0] & \text{on} \ J \times \Sigma, \\
v(t) &= e^{\Delta \Omega t} (v_0|_{\partial \Omega}) & \text{on} \ J \times \partial \Omega, \\
v(0) &= v_0 & \text{in} \ \Omega^t,
\end{align*}
\]

where \( \Delta_\Sigma \), respectively on \( \Delta_{\partial \Omega} \) denotes the Laplace-Beltrami operator on \( \Sigma \) respectively \( \partial \Omega \) admits a unique solution \( v \in H_p^1(J; L_p(\Omega)^n) \cap L_p(J; H_p^{1-2/p}(\Omega^t)) \). We may now define

\[ f := \text{div} v \in H_p^1(J; \dot{H}_p^{-1}(\Omega)) \cap L_p(J; H_p^1(\Omega^t)) \]

and by construction we have

\[ f(0) = \text{div} v_0 = \Delta \phi_0 = f_0. \]

Finally, we define \( g \) by means of \( g = e^{\Delta \Omega t} g_0 \) to the result

\[ g \in W_p^{1/2-1/2p}(J; L_p(\Sigma)) \cap L_p(J; W_p^{1-1/p}(\Sigma)) \]

by the properties of the analytic semigroup \( e^{\Delta \Sigma t} \).
Next we introduce the linear operator $L = (L_1, \ldots, L_4)$ defined by the left-hand side of (2.2), i.e.

\begin{align*}
L_1(u, \pi) &:= \rho \partial_t u - \mu \Delta u + \nabla \pi, \\
L_2(u) &:= (\nabla |u|), \\
L_3(u, q, h) &:= [-\mu(\nabla u + [\nabla u]^T)]\nu_\Sigma + q \nu_\Sigma - (\Delta \Sigma) h \nu_\Sigma, \\
L_4(u, h) &:= \partial_t h - (u|\nu_\Sigma) + (b|\nabla \Sigma)h,
\end{align*}

and the nonlinearity $N = (N_1, \ldots, N_4)$ defined by the right hand side of (2.2), i.e.

\begin{align*}
N_1(u, \pi, h) &:= F(h, u) \nabla u + M_4(h) : \nabla^2 u + M_1(h) \nabla \pi, \\
N_2(u, h) &:= M_1(h) : \nabla u, \\
N_3(u, h) &:= G_z(h) \nabla u + (G_p(h) \nabla u + G_z(h)) \nu_\Sigma, \\
N_4(u, h) &:= ([M_0(h) - I] \nabla \Sigma |u|) + (u - b|\nabla \Sigma h).
\end{align*}

For $J = [0, a]$ let the solution spaces be defined by

\begin{align*}
E_1(a) &:= \{ u \in H^2_p(J; L_p(\Omega)^n) \cap L_p(J; H^2_p(\Omega \setminus \Sigma)^n) : u = 0 \text{ on } \partial \Omega, \quad [u] = 0 \}, \\
E_2(a) &:= L_p(J; H^1_p(\Omega \setminus \Sigma)), \\
E_3(a) &:= W^{1/2-1/p}_p(J; L_p(\Sigma)^n) \cap L_p(J; W^{1-1/p}_p(\Sigma)^n), \\
E_4(a) &:= W^{2-1/2p}_p(J; L_p(\Sigma)) \cap H^1_p(J; W^{2-1/p}_p(\Sigma)) \cap L_p(J; W^{3-1/p}_p(\Sigma)).
\end{align*}

We abbreviate

\begin{align*}
E(a) &:= \{ (u, \pi, q, h) \in E_1(a) \times E_2(a) \times E_3(a) \times E_4(a) : [\pi] = q \},
\end{align*}

and equip $E_1(a)$, $E_2(a)$, $E_3(a)$ and $E_4(a)$ with their natural norms, which turn them into Banach spaces; $E(a)$ carries the natural norm of the underlying product space. A left subscript 0 always means that the time trace of the function is zero whenever it exists. Furthermore, the data spaces are defined by

\begin{align*}
F_1(a) &:= L_p(J; L_p(\Omega)^n), \\
F_2(a) &:= H^1_p(J; H^{-1}_p(\Omega)) \cap L_p(J; H^1_p(\Omega \setminus \Sigma)), \\
F_3(a) &:= W^{1/2-1/2p}_p(J; L_p(\Sigma)^n) \cap L_p(J; W^{1-1/p}_p(\Sigma)^n), \\
F_4(a) &:= W^{1-1/2p}_p(J; L_p(\Sigma)) \cap L_p(J; W^{2-1/p}_p(\Sigma))
\end{align*}

and

\begin{align*}
F(a) &:= F_1(a) \times F_2(a) \times F_3(a) \times F_4(a),
\end{align*}

where we equip these spaces again with their natural norms. The generic elements of $F(a)$ are $(f, f_4, g, g_h)$.

To shorten the notation we set $z = (u, \pi, q, h) \in E(a)$ and reformulate the quasi-linear problem (2.2) as

\begin{equation}
Lz = N(z), \quad (u(0), h(0)) = (u_0, h_0).
\end{equation}

From Section 3 we already know that $L : E(a) \to F(a)$ is bounded and linear and that $L : 0E(a) \to 0F(a)$ is an isomorphism, for each $a > 0$.

Concerning the nonlinearity $N$, the following result has been shown in [20, Proposition 4.1] for the case where $\Sigma$ is a graph over $\mathbb{R}^n$. 

\begin{align*}
\text{(4.2)}
\end{align*}
Proposition 4.2. Suppose $p > n + 2$ and $b \in F_4(a)^n$. Then

$$N \in C^\infty(\mathbb{E}(a), F(a)), \quad a > 0. \quad (4.3)$$

Let $DN(u, \pi, q, h)$ denote the Fréchet derivative of $N$ at $(u, \pi, q, h) \in \mathbb{E}(a)$. Then $DN(u, \pi, q, h) \in \mathcal{L}(\mathbb{E}(a), _0F(a))$, and for any number $a_0 > 0$ there is a positive constant $M_0 = M_0(a_0, p)$ such that

$$|DN(u, \pi, q, h)|_{\mathcal{L}(\mathbb{E}(a), _0F(a))} \leq M_0\left[|b - u|_{BC(J, BC)} + |(u, \pi, q, h)|_{\mathbb{E}(a)}\right]$$

$$+ M_0\left[(\nabla h)_{BC(J, BC)} + |h|_{\mathbb{E}(a)} + |u|_{BC(J, BC)}\right]|u|_{\mathbb{E}(a)}$$

$$+ M_0\left[P(|\nabla h|_{BC(J, BC)}\right)\nabla h|_{BC(J, BC)} + Q(|\nabla h|_{BC(J, BC)}|h|_{\mathbb{E}(a)})|h|_{\mathbb{E}(a)}\right]$$

for all $(u, \pi, q, h) \in \mathbb{E}(a)$ and all $a \in (0, a_0]$. Here, $P$ and $Q$ are fixed polynomials with coefficients equal to one.

The proof carries over to the general case considered here. The basic ingredients are still the polynomial structure of the nonlinearity $N$ w.r.t. $u$ and $\pi$, which is the same as in [20], and the embeddings

$$W_p^{n+2} = BUC^{n+2}(\Omega \setminus \Sigma), \quad W_p^{n+2} = BUC^{n+2}(\Omega \setminus \Sigma),$$

with $\alpha = 1 - (n + 2)/p > 0$, which show that $\mathbb{E}(a)$ and $F_4(a)$ are Banach algebras. The difference lies only in the operators $M_j(h)$ which are more complicated in the case of general domains, but analytic in $h$.

Now, we are able to establish local well-posedness.

Theorem 4.3. Fix $p > n + 2$, let $\partial \Omega \in C^3$, and suppose

$$\Gamma_0 \in W_p^{n-2/p}, \quad u_0 \in W_p^{2-2/p}(\Omega \setminus \Gamma_0)^n.$$ 

Assume the compatibility conditions

$$\text{div } u_0 = 0 \quad \text{in } \Omega \setminus \Gamma_0, \quad u_0 = 0 \quad \text{on } \partial \Omega,$$

$$[P_{\Gamma, \alpha} \mu (\nabla u_0 + [\nabla u_0]^T)v_{\Gamma_0}] = 0, \quad [u_0] = 0 \quad \text{on } \Gamma_0,$$

where $P_{\Gamma, \alpha} = I - v_{\Gamma_0} \times v_{\Gamma_0}$.

Then there exists $a = a(u_0, \Gamma_0) > 0$ and a unique classical solution $(u, \pi, \Gamma)$ of (4.1) on $(0, a)$. The set

$$\Upsilon = \bigcup_{t \in (0, a)} \{t\} \times \Gamma(t)$$

is a real analytic manifold, and with

$$\bar{U} \setminus \{t \times \Omega \times x \notin \Gamma(t)\},$$

the function $(u, \pi) : \bar{U} \rightarrow \mathbb{R}^{n+1}$ is real analytic. The transformed solution $(\bar{u}, \bar{\pi}, \bar{q}, h)$ belongs to the space $\mathbb{E}(a)$.

Proof. We consider the transformed problem.

Step 1. Let $h_0 \in W_p^{n+2} = BUC^{n+2}(\Omega \setminus \Sigma)$ and $u_0 \in W_p^{2-2/p}(\Omega \setminus \Sigma)^n$ be given, such that the compatibility conditions are satisfied, and $|h_0|_{BUC^{n+2}(\Sigma)} \leq \eta$. Let $J_0 = [0, a_0]$ and

$$f^* \in W_p^{1/2-1/2p}(J_0; H_p^{1/2-1/2p}(\Omega \setminus \Sigma)) \cap L_p(J_0; H_p^{1/2}(\Omega \setminus \Sigma)),$$

$$g^* \in W_p^{1/2-1/2p}(J_0; L_p(\Sigma)) \cap L_p(J_0; W_p^{1/2-1/p}(\Sigma)).$$
be extensions of
\[(\nabla|u_0) \in H^{-1}_p(\Omega) \cap W^{1-2/p}(\Omega \backslash \Sigma) \quad \text{and} \quad \nabla[\mu(\nabla u_0 + [\nabla u_0]^T)\nu_2] \in W^{1-3/p}(\Sigma),\]
which exist due to Proposition 4.1. Further choose an extension \( \tilde{u} \in \mathbb{E}_1(a_0) \) of \( u_0 \) and set \( b = \tilde{u} \) restricted to \([0, a_0] \times \Sigma\). With these extensions we may solve the problem
\[Lz^* = (0, f^*_d, g^*, 0), \quad (u^*(0), h^*(0)) = (u_0, h_0),\]
since all regularity and compatibility conditions of Theorem 4.3 are satisfied.

**Step 2.** We rewrite problem \((2.2)\) as
\[Lz = N(z + z^*) - Lz^* =: K(z), \quad z \in \mathbb{E}(a)\]
and observe, that the solution is given as \( z = L^{-1}K(z) \), since Theorem 4.3 implies that \( L : \mathbb{E}(a) \to \mathbb{F}(a) \) is an isomorphism with
\[|L^{-1}|_{\mathbb{E}(a)} \leq M, \quad a \in (0, a_0],\]
where \( M \) is independent of \( a \). Thanks to Proposition 4.2 and due to \( K(0) = N(z^*) - Lz^* \), we may choose \( a \in (0, a_0] \) and \( r > 0 \) sufficiently small such that
\[|K(0)|_{\mathbb{F}(a)} \leq \frac{r}{2M}, \quad |DK(z)|_{\mathbb{E}(a), \mathbb{F}(a)} \leq \frac{1}{2M}, \quad z \in \mathbb{E}(a), \quad |z|_{\mathbb{E}(a)} \leq r,\]
hence
\[|K(z)|_{\mathbb{F}(a)} \leq \frac{r}{M},\]
which ensures, that \( L^{-1}K : \mathbb{E}_p(0) \to \mathbb{E}_p(0) \) is a strict contraction; see also [20]. Thus, we may employ the contraction mapping principle to obtain a unique solution on the time interval \([0, a]\).

**Step 3.** By Proposition 4.2 the right-hand side \( N \) is real analytic, and hence, we obtain analyticity of \((u, \pi, q, h)\) in space and time by the parameter trick as shown in [13] Theorem 6.3; cf. also [13, 14]. □

At the end of this section, we want to mention an extension of Theorem 4.3 to the case of time-weighted \( L_p \)-spaces. For this purpose, for \( \mu \in (1/p, 1] \) we define \( \mathbb{E}_\mu(a) \) by means of
\[z \in \mathbb{E}_\mu(a) \quad \Leftrightarrow \quad t^{1-\mu}z \in \mathbb{E}(a),\]
and similarly we define \( \mathbb{F}_\mu(a) \). Thus \( \mathbb{E}_1(a) = \mathbb{E}(a), \mathbb{F}_1(a) = \mathbb{F}(a) \). Such time-weights are useful to relax the regularity of the initial values, but maintaining the regularity of the solution for \( t \geq \delta > 0 \) for arbitrary small positive \( \delta \). More precisely, we have the following result.

**Corollary 4.4.** Fix \( p > n + 2, \mu \in \left( \frac{1}{p} + \frac{n+2}{2p}, 1 \right), \) let \( \partial \Omega \in C^3 \), and suppose
\[\Gamma_0 \in W^{2+\mu-2/p}_p, \quad u_0 \in W^{2\mu-2/p}(\Omega \backslash \Gamma_0)\]
are subject to the compatibility conditions \((4.1)\).

Then there exists a \( a = a(u_0, \Gamma_0) > 0 \) and a unique solution \((\bar{u}, \bar{\pi}, \bar{q}, h) \in \mathbb{E}_\mu(a) \) of the transformed problem \((2.2)\), which depends continuously on the initial data.
This result is proved in the same way as Theorem 4.3, taking into account that the restriction $\mu > \frac{1}{2} + \frac{n-2}{2p}$ still ensures the embeddings

$$W^{2\mu - 2/p}_p(\Omega, \Sigma) \hookrightarrow \text{BUC}^1(\Omega \setminus \Sigma), \quad W^{2+\mu - 2/p}_p(\Sigma) \hookrightarrow \text{BUC}^2(\Sigma),$$

which are crucial for Proposition 4.2 to be valid also in the corresponding time-weighted spaces. Further, the result about the linear problem remains valid in these time-weighted spaces as well. This can be shown in the same way as in the case $\mu = 1$, taking into account that the operator $d/dt$ admits an $H^\infty$-calculus in $L_{p,\mu}(J; X)$ with angle $\pi/2$, provided $\mu > 1/p$ and $X$ is UMD-Banach space; cf. [17]. We refrain from giving more details, here; see also [14].

5. Semiflow, Energy Functional and Equilibria

5.1. The Induced Semiflow. Recall that the closed $C^2$-hyper-surfaces contained in $\Omega$ form a $C^2$-manifold, which we denote by $\mathcal{MH}^2(\Omega)$. The metric on $\mathcal{MH}^2(\Omega)$ is defined by

$$d(\Sigma_1, \Sigma_2) := d_H(\mathcal{N}^2\Sigma_1, \mathcal{N}^2\Sigma_2), \quad \Sigma_1, \Sigma_2 \in \mathcal{MH}^2(\Omega).$$

The charts are the parameterizations over a given real analytic hyper-surface $\Sigma$, as described in Section 2, and the tangent space at $\Sigma$ consist of the normal vector fields on $\Sigma$ of class $C^2$. This way $\mathcal{MH}^2(\Omega)$ becomes a Banach manifold.

Let $d_\Sigma(x)$ denote the signed distance for $\Sigma$ as introduced in Section 2. We may then define a level function $\varphi_\Sigma$ by means of

$$\varphi_\Sigma(x) = g(d_\Sigma(x)), \quad x \in \mathbb{R}^n,$$

where

$$g(s) = s(1 - \chi(s/a)) + \chi(s/a)\text{sgn }s, \quad s \in \mathbb{R},$$

and $\chi$ denotes the cut-off function defined in Section 2. Then it is easy to see that $\Sigma = \varphi_\Sigma^{-1}(0)$, and $\nabla \varphi_\Sigma(x) = \nu_\Sigma(x)$, for each $x \in \Sigma$. Moreover, $\mu = 0$ is an eigenvalue of $\nabla^2 \varphi_\Sigma(x)$ with eigenfunction $\nu_\Sigma(x)$, the remaining eigenvalues of $\nabla^2 \varphi_\Sigma(x)$ are the principal curvatures $\kappa_j(x)$ of $\Sigma$ at $x \in \Sigma$.

Consider the subset $\mathcal{MH}^2(\Omega, r)$ of $\mathcal{MH}^2(\Omega)$ which consists of all $\Gamma \in \mathcal{MH}^2(\Omega)$ such that $\Gamma \subset \Omega$ satisfies the ball condition with fixed radius $r > 0$. This implies in particular that $\text{dist}(\Gamma, \partial \Omega) \geq r$ and all principal curvatures of $\Gamma \in \mathcal{MH}^2(\Omega, r)$ are bounded by $r$. Further, the level functions $\varphi_\Gamma = g \circ d_\Gamma$ are well defined for $\Gamma \in \mathcal{MH}^2(\Omega, r)$, and form a bounded subset of $C^2(\Omega)$. The map $\Phi : \mathcal{MH}^2(\Omega, r) \to C^2(\Omega)$ defined by $\Phi(\Gamma) = \varphi_\Gamma$ is an isomorphism of the metric space $\mathcal{MH}^2(\Omega, r)$ onto $\Phi(\mathcal{MH}^2(\Omega, r)) \subset C^2(\Omega)$.

Let $s - (n-1)/p > 2$; for $\Gamma \in \mathcal{MH}^2(\Omega, r)$, we define $\Gamma \in W_p^s(\Omega, r)$ if $\varphi_\Gamma \in W_p^s(\Omega)$. In this case the local charts for $\Gamma$ can be chosen of class $W_p^s$ as well. A subset $A \subset W_p^s(\Omega, r)$ is said to be (relatively) compact, if $\Phi(A) \subset W_p^s(\Omega)$ is (relatively) compact. Finally, we define

$$\text{dist}_{W_p^s}(\Gamma, \Sigma) := |\varphi_\Gamma - \varphi_\Sigma|_{W_p^s(\Omega)}$$

for $\Gamma, \Sigma \in \mathcal{MH}^2(\Omega, r)$.

As an ambient space for the phase-manifold $\mathcal{P}M$ of the two-phase Navier-Stokes problem with surface tension we consider the product space $C(\Omega)^n \times \mathcal{MH}^2(\Omega)$.
We define \( \mathcal{PM} \) as follows.

\[
\mathcal{PM} := \{(u, \Gamma) \in C(\bar{\Omega}) \times \mathcal{M}H^2(\Omega) : u \in W^{2-2/p}_p(G\setminus\Gamma)^n, \Gamma \in W^{3-2/p}_p, \forall \nu \in \partial\Omega, \mathcal{P}_\Gamma[\mu(\nabla u + [\nabla u]^T)]\nu_T = 0 \text{ on } \Gamma \}. \tag{5.1}
\]

The charts for this manifold are obtained by the charts induced by \( \mathcal{M}H^2(\Omega) \), followed by a Hanzawa transformation; see Section 2.

Observe that the energy functional is a Ljapunov functional, in fact, even a strict one. We have the following result.

Let \( p > n + 2 \). Then the two-phase Navier-Stokes problem with surface tension generates a local semiflow on the phase-manifold \( \mathcal{PM} \). Each solution \( (u, \Gamma) \) exists on a maximal time interval \( [0, t_*) \).

### 5.2. The Pressure

The pressure does not occur explicitly as a variable in the local semiflow, the latter is only formulated in terms of the velocity field \( u \) and the free boundary \( \Gamma \). Actually, at every instant \( t \) the pressure \( \pi \) can be reconstructed from the semiflow. In fact, fix any \( t \in (0, t_*) \) and consider \( \phi \in H^1_p(\Omega) \). Then we have by the divergence theorem

\[
(u(t)\nabla \phi)_{L^2(\Omega)} = -\langle \text{div } u(t)\phi \rangle_{L^2(\Omega)} - \langle [u(t)\nu]\phi \rangle_{L^2(\Gamma)} = 0,
\]

hence also \( \partial_t u \nabla \phi = 0 \). This implies, multiplying the momentum balance divided by \( \rho \) with \( \nabla \phi \) in \( L^2(\Omega) \)

\[
\left( \nabla \frac{\pi}{\rho} \bigg| \nabla \phi \right)_{L^2(\Omega)} = \left( \frac{\mu}{\rho} \Delta u - u \cdot \nabla u \bigg| \nabla \phi \right)_{L^2(\Omega)}.
\]

On the other hand, multiplying the stress boundary condition by \( \nu_T \) yields

\[
[\pi] = \sigma H_T + ([\mu(\nabla u + [\nabla u]^T)]\nu_T|\nu_T) \]

on \( \Gamma \). Thus \( \pi \) must satisfy the following problem.

\[
\left( \nabla \frac{\pi}{\rho} \bigg| \nabla \phi \right)_{L^2(\Omega)} = \left( \frac{\mu}{\rho} \Delta u - u \cdot \nabla u \bigg| \nabla \phi \right)_{L^2(\Omega)}, \quad \phi \in H^1_p(\Omega) \tag{5.2}
\]

\[
[\pi] = \sigma H_T + ([\mu(\nabla u + [\nabla u]^T)]\nu_T|\nu_T), \quad x \in \Gamma.
\]

Theorem 8.3 implies that this problem has a unique solution \( \pi \in \dot{H}^1_p(\Omega\setminus\Gamma) \). Thus the pressure is uniquely defined (up to a constant) by the semiflow and can be obtained by solving the transmission problem (5.2).

### 5.3. The Energy Functional

Define the energy functional by means of

\[
\Phi(u, \Gamma) := \frac{1}{2} \rho^{1/2} u_{L^2(\Omega)}^2 + \sigma|\Gamma(t)|.
\]

Then

\[
\partial_t \Phi(u, \Gamma) + 2\mu^{1/2}E_{L^2(\Omega)}^2 = 0,
\]

hence the energy functional is a Ljapunov functional, in fact, even a strict one. We have the following result.
Proposition 5.2. Let \( \rho_i, \mu_i, \sigma > 0 \) be constants. Then

(a) The energy equality is valid for smooth solutions.
(b) The equilibria are zero velocities, constant pressures in the components of the phases, the dispersed phase is a union of nonintersecting open balls.
(c) The energy functional is a strict Lyapunov-functional.
(d) The critical points of the energy functional for constant phase volumes are precisely the equilibria.

This result is a special case of [3, Theorem 3.1].

Remark 5.3. (i) Let us point out that in equilibrium the dispersed phase consists of at most countably many disjoint balls \( B_{R_i}(x_i) \). If there are infinitely many of them, then \( R_i \to 0 \) as \( i \to \infty \), hence the corresponding curvatures \( H_i = -(n-1)/R_i \) tend to infinity, as well as the pressures inside these balls. This is due to the model assumption that there is no phase transition. On the other hand, phase transition will occur at very high pressure levels. To avoid this contradiction, in the sequel we consider only equilibria in which the dispersed phase consists of only finitely many balls. Note also that the free boundary will not be of class \( C^2 \) if \( \Omega_1 \) has infinitely many components.

(ii) There is another pathological case which we exclude in the sequel, namely if the dispersed phase contains balls touching each other. This can only happen if the radii of these balls are equal, otherwise the pressure jump would not be constant on \( \Gamma \). Physically one would expect such an equilibrium to be unstable, but at present we are not able to handle this case. Observe that also in such a situation the free boundary \( \Gamma \) is not a manifold of class \( C^2 \).

(iii) Of course neither (i) or (ii) occurs if we assume that \( \Omega_1 \) is connected, the continuous phase enjoys this property anyway.

6. The Stability Result

Assuming, for simplicity, that the phases are connected, we denote by

\[ E := \{ (0, S_R(x_0)) : x_0 \in \Omega, R > 0, \bar{B}_R(x_0) \subset \Omega \} \]

the set of equilibria without boundary contact. Note that \( E \) forms a real analytic manifold of dimension \( n + 1 \). Here \( n \) dimensions come from the coordinates of the center \( x_0 \) and one from the radius \( R \) of the sphere \( S_R(x_0) \).

Fix any such equilibrium \((0, \Sigma) \in E\). We consider the behaviour of the solutions near this steady state. Suppose \( p > n + 2 \), let \( \partial \Omega \in C^3 \), and consider initial data \((u_0, \Gamma_0) \in PM\).

Here we have to use the full linearization of the problem at an equilibrium \((0, \Sigma)\) i.e. at \((u, h) = (0, 0)\), and for this reason we have to replace \( \Delta_\Sigma \) in the linear problem [3,1] by

\[ A_\Sigma = H'_\Gamma(0) = \frac{n-1}{R^2} + \Delta_\Sigma. \]
This results in the problem

\begin{align*}
\rho \partial_t v - \mu \Delta v + \nabla q &= \rho f_v \quad \text{in } \Omega \setminus \Sigma, \\
\text{div } v &= f_d \quad \text{in } \Omega \setminus \Sigma, \\
-\mu (|\nabla v| + |\nabla v|^T) + q)\nu_\Sigma - \sigma (A_\Sigma h)\nu_\Sigma &= g \quad \text{on } \Sigma, \\
[v] &= 0 \quad \text{on } \Sigma, \\
v &= 0 \quad \text{on } \partial \Omega, \\
\partial_t h - (v|\nu_\Sigma) &= g_h \quad \text{on } \Sigma, \\
v(0) &= v_0 \quad \text{in } \Omega \setminus \Sigma, \\
h(0) &= h_0 \quad \text{on } \Sigma.
\end{align*}

(6.1)

It is well-known that $A_\Sigma$ is selfadjoint, negative semidefinite on functions with zero mean, and has compact resolvent in $L^2(\Sigma)$. $\lambda_0 = 0$ is an eigenvalue with eigenspace of dimension $n$, spanned by the spherical harmonics of degree one. $\lambda_{-1} = (n-1)/R^2$ is also an eigenvalue, its eigenspace is one-dimensional and consists of the constants.

As a base space for our analysis we use

$$X_0 = L_{p,\sigma}(\Omega)^n \times W^{2-1/p}(\Sigma),$$

where the subscript $\sigma$ means solenoidal, and we set

$$\bar{X}_1 = H^2_p(\Omega \setminus \Sigma)^n \times W^{3-1/p}(\Sigma).$$

Define a closed linear operator in $X_0$ by means of

$$A(v,h) = \left(-\left(\frac{\mu}{\rho}\right)\Delta v + \rho^{-1}\nabla q, -v|\nu_\Sigma\right),$$

with domain $X_1 := D(A) \subset \bar{X}_1$ defined by

$$D(A) = \{(v,h) \in \bar{X}_1 \cap X_0 : v = 0 \text{ on } \partial \Omega, \ [v] = 0 \text{ and}$$

$$[P_\Sigma \mu (\nabla v + \nabla v^T)\nu_\Sigma] = 0 \text{ on } \Sigma\},$$

where as before $P_\Sigma$ means the projection onto the tangent space of $\Sigma$. Here $q \in H^1_p(\Omega \setminus \Sigma)$ is determined as the solution of the transmission problem

$$\left(\frac{\nabla q}{\rho} , \nabla \phi\right)_{L^2} = \left(\frac{\mu}{\rho}\Delta v, \nabla \phi\right)_{L^2}, \quad \phi \in W^1_p(\Omega),$$

$$[q] = [\mu((\nabla v + |\nabla v|^T)\nu_\Sigma)] + \sigma A_\Sigma h \quad \text{on } \Sigma,$$

which is well-defined (up to a constant) by Theorem 8.5. This implies

$$\frac{1}{\rho} \nabla q = T_1 \left(\frac{\mu}{\rho}\Delta v\right) + T_2 \left(\|\mu(\nabla v + |\nabla v|^T)\nu_\Sigma\| + \sigma A_\Sigma h\right).$$

Then with $z = (v,h)$ and $f = (f_v, g_h)$ as well as $z_0 = (v_0, h_0)$, system (6.1) can be rewritten as the abstract evolution equation

$$\dot{z} + Az = f, \quad t > 0, \quad z(0) = z_0,$$

(6.2)

in $X_0$, provided $f_d \equiv g \equiv 0$; see also 25.

Since by Theorem 3.1 problem (6.1) has maximal $L_p$-regularity, the abstract problem (6.2) has maximal $L_p$-regularity, as well. In particular, $-A$ generates an analytic $C_0$-semigroup in $X_0$; see e.g. 16, Proposition 1.2. In addition we have the following result.
Proposition 6.1. Let $\rho_i, \mu_i > 0, \sigma > 0$ be constants, $p \in (1, \infty)$, and let $X_0, A, X_1 := D(A)$ be defined as above. Then the following holds.
(a) The linear operator $-A$ generates a compact analytic $C_0$-semigroup in $X_0$ which has the property of maximal $L_p$-regularity.
(b) The spectrum of $A$ consists of countably many eigenvalues with finite algebraic multiplicity and is independent of $p$.
(c) $A$ has no eigenvalues $\lambda$ with nonnegative real part other than $\lambda = 0$.
(d) $\lambda = 0$ is a semisimple eigenvalue with multiplicity $n + 1$.
(e) The eigenspace $N(A)$ is isomorphic to the tangent space $T_z E$ of $E$ at the given equilibrium $z^* = (0, \Sigma)$, where $\Sigma = S_R(x_0)$.
(f) The restriction of $e^{-At}$ to $R(A)$ is exponentially stable.

This result is a special case of [11, Theorem 4.1]. It shows that equilibria $(0, \Sigma) \in E$ are normally stable, hence allows for the use of the generalized principle of linearized stability to obtain our main result on stability and convergence.

The following result concerns the stationary Stokes problem
\begin{equation}
\begin{aligned}
\rho \omega u - \mu \Delta u + \nabla \pi &= 0 \quad \text{in} \quad \Omega \setminus \Sigma, \\
\mathrm{div} \ u &= f_d \quad \text{in} \quad \Omega \setminus \Sigma, \\
-\mathcal{P}_{\Sigma} [\mu (\nabla u + [\nabla u]^T)] \nu_{\Sigma} &= g_\tau \quad \text{on} \quad \Sigma, \\
-([\mu (\nabla u + [\nabla u]^T)] \nu_{\Sigma} + [\pi]) &= g_\nu \quad \text{on} \quad \Sigma, \\
[u] &= 0 \quad \text{on} \quad \Sigma, \\
u &= 0 \quad \text{on} \quad \partial \Omega.
\end{aligned}
\end{equation}

(6.3)

It is needed in the proof of the main result of this section.

Proposition 6.2. Let $p > 3$ be fixed, and assume that $\rho_i$ and $\mu_i$ are positive constants for $i = 1, 2$, and that $\omega > 0$ is large enough. Then the stationary Stokes problem with free boundary (6.3) admits a unique solution $(u, \pi, [\pi])$ with regularity
\begin{equation}
u \in W_p^{2-2/p}(\Omega \setminus \Sigma)^n, \quad \pi \in W_p^{1-2/p}(\Omega \setminus \Sigma), \quad [\pi] \in W_p^{1-3/p}(\Sigma),
\end{equation}

(6.4)

if and only if the data $(f_d, g)$ satisfy the following regularity conditions:
(a) $f_d \in W_p^{1-2/p}(\Omega \setminus \Sigma) \cap \dot{H}_p^{-1}(\Omega)$,
(b) $g = (g_\tau, g_\nu) \in W_p^{1-3/p}(\Sigma)^n$.

The solution map $[(f_d, g) \mapsto (u, \pi, [\pi])]$ is continuous between the corresponding spaces.

The proof of this elliptic problem is similar to that of Theorem 3.1. We refer to [21] for the case of a flat interface.

The main result of this section is the following.

Theorem 6.3. The equilibrium $(0, \Sigma) \in E$ is stable in the sense that for each $\epsilon \in (0, \epsilon_0]$ there exists $\delta(\epsilon) > 0$ such that for all initial values $(u_0, \Gamma_0)$ subject to
\begin{equation}
\begin{aligned}
\text{dist}_{W_p^{2-2/p}(\Gamma_0, \Sigma)} \leq \delta(\epsilon) \quad \text{and} \quad \|u_0\|_{W_p^{2-2/p}(\Omega, \Gamma_0)} \leq \delta(\epsilon)
\end{aligned}
\end{equation}

there exists a unique global solution $(u(t), \Gamma(t))$ of the problem, and it satisfies
\begin{equation}
\begin{aligned}
\text{dist}_{W_p^{2-2/p}(\Gamma(t), \Sigma)} \leq \epsilon \quad \text{and} \quad \|u(t)\|_{W_p^{2-2/p}(\Omega, \Gamma(t))} \leq \epsilon, \ t \geq 0.
\end{aligned}
\end{equation}
Moreover, as \( t \to \infty \) the solutions \((u(t), \Gamma(t))\) converges to an equilibrium \((0, \Sigma_{\infty})\) in the same topology, i.e.

\[
\lim_{t \to \infty} \text{dist}_{W^{3-2/p}_p(\Gamma(t), \Sigma_{\infty})} + \|u(t)\|_{W^{2-2/p}_p(\Omega \setminus \Gamma(t))} = 0.
\]

The convergence is at exponential rate.

**Proof.** 1. To prepare, as in [21] we first parameterize the nonlinear phase manifold locally near \((0, \Sigma)\) over

\[
X_\gamma := \{(u, h) \in [W^{2-2/p}_p(\Omega \setminus \Sigma) \times W^{3-2/p}_p(\Sigma)] \cap X_0 : u = 0 \text{ on } \partial \Omega, \quad [u] = 0, \quad \mathcal{P}_2[\mu(\nabla u + [\nabla u]^T)]|_{\nu_2} = 0 \text{ on } \Sigma, \quad \bar{\omega} > 0 \text{ on } \Sigma \}
\]

In particular, this will show that \( X_\gamma \) is isomorphic to the tangent space \( T_{(0, \Sigma)} \mathcal{P}M \).

For this purpose fix \( \omega > 0 \) and solve for given \( \bar{\omega} := (\bar{u}, \bar{h}) \in B_r^{X_\gamma}(0) \) the problem

\[
\begin{align*}
\rho \omega \ddot{u} - \mu \Delta \ddot{u} + \nabla \bar{\pi} &= 0 & \text{in } \Omega \setminus \Sigma, \\
\text{div } \ddot{u} &= M_1(h) : \nabla \ddot{u} & \text{in } \Omega \setminus \Sigma, \\
-\mathcal{P}_2[\mu(\nabla \bar{u} + [\nabla \bar{u}]^T)]|_{\nu_2} &= G_\tau(h) \nabla \bar{u} & \text{on } \Sigma, \\
-([\mu(\nabla \bar{u} + [\nabla \bar{u}]^T)]|_{\nu_2} + [\bar{\pi}]) &= G_\nu(h) \nabla \bar{u} + G_\gamma(h) & \text{on } \Sigma, \\
[u] &= 0 & \text{on } \Sigma, \\
\bar{u} &= 0 & \text{on } \partial \Omega,
\end{align*}
\]

where \( u = \bar{u} + \hat{u} \) and \( h = \hat{h} \), i.e. \( \bar{h} = 0 \). We write this equation in short hand notation

\[
\begin{align*}
\mathfrak{L}_\omega \ddot{u} &= \mathfrak{N}(\bar{\omega} + \bar{\dot{\omega}}) \text{ in } \mathcal{X}_\gamma, \\
\text{div } \ddot{u} &= M_1(h) : \nabla \ddot{u} & \text{in } \mathcal{X}_\gamma, \\
-\mathcal{P}_2[\mu(\nabla \bar{u} + [\nabla \bar{u}]^T)]|_{\nu_2} &= G_\tau(h) \nabla \bar{u} & \text{on } \mathcal{X}_\gamma, \\
-([\mu(\nabla \bar{u} + [\nabla \bar{u}]^T)]|_{\nu_2} + [\bar{\pi}]) &= G_\nu(h) \nabla \bar{u} + G_\gamma(h) & \text{on } \mathcal{X}_\gamma, \\
[u] &= 0 & \text{on } \mathcal{X}_\gamma, \\
\bar{u} &= 0 & \text{on } \partial \mathcal{X}_\gamma,
\end{align*}
\]

where \( \mathfrak{L}_\omega \) is real analytic and \( \mathfrak{N}(0) = 0 \); see Section 4. \( \mathfrak{L}_\omega \) is invertible by Proposition [6,2] hence the implicit function theorem yields a unique solution \( \bar{u} = \phi(\bar{u}, \hat{h}) \in \mathcal{X}_\gamma \) near 0. \( \phi \) is real analytic as well, and satisfies \( \phi'(0) = 0 \). Then we define

\[
\Phi(\bar{u}, \hat{h}) = (\bar{u}, \hat{h}) + (\phi(\bar{u}, \hat{h}), 0).
\]

Obviously, \( \Phi \) is real analytic, \( \Phi'(0) = I, \Phi(B_r^{X_\gamma}(0)) \subset \mathcal{P}M \), and \( \Phi \) is injective.

Hence it remains to show local surjectivity near 0. Suppose that \( \bar{\omega} := (\bar{u}, \hat{h}) \in \mathcal{P}M \) has sufficiently small norm. Solving the problem

\[
\begin{align*}
\rho \omega u - \mu \Delta u + \nabla \pi &= 0 & \text{in } \Omega \setminus \Sigma, \\
\text{div } u &= M_1(h) : \nabla u & \text{in } \Omega \setminus \Sigma, \\
-\mathcal{P}_2[\mu(\nabla u + [\nabla u]^T)]|_{\nu_2} &= G_\tau(h) \nabla u & \text{on } \Sigma, \\
-([\mu(\nabla u + [\nabla u]^T)]|_{\nu_2} + [\pi]) &= G_\nu(h) \nabla u + G_\gamma(h) & \text{on } \Sigma, \\
[u] &= 0 & \text{on } \Sigma, \\
u &= 0 & \text{on } \partial \Omega,
\end{align*}
\]

by means of Proposition [6,2] \( (\bar{u}, \hat{h}) := (\bar{u} - u, \hat{h}) \) belongs to \( X_\gamma \) and \( \phi(\bar{u}, \hat{h}) = u \), showing surjectivity of \( \Phi \) near 0. In particular, \( \mathcal{P}M \) is a real analytic manifold near \((0, \Sigma)\).
2. Let \((u, \pi, h)\) be a solution on its maximal time interval \([0, t_\ast)\). In this step we decompose \(u = \bar{u} + \check{u}, \pi = \bar{\pi} + \check{\pi}, h = \bar{h} + \check{h}\), where \((\bar{u}, \bar{\pi}, \bar{h})\) solves the problem

\[
\begin{align*}
\rho \omega \bar{u} + \rho \partial_t \bar{u} - \mu \Delta \bar{u} + \nabla \bar{\pi} &= F_u(u, \pi, h) & \text{in } \Omega \setminus \Sigma, t > 0, \\
\text{div } \bar{u} &= M_1(h) : \nabla u & \text{in } \Omega \setminus \Sigma, t > 0, \\
-P_{\Sigma} [\mu E(\bar{u})]\nu_{\Sigma} &= G_\gamma(h) \nabla u & \text{on } \Sigma, t > 0,
\end{align*}
\]

\(-([\mu E(\bar{u})]\nu_{\Sigma}) + [\pi] - \sigma A_{\Sigma} \bar{h} = G_\gamma(h) \nabla u\) on \(\Sigma, t > 0\),

\[
[\bar{u}] = 0 \quad \text{on } \Sigma, t > 0,
\]

\[
\bar{u} = 0 \quad \text{on } \partial \Omega, t > 0,
\]

\[
\omega \check{h} + \partial_t \check{h} - (\check{u})\nu_{\Sigma} = (M_0(h) \nabla \check{h}|u) \quad \text{on } \Sigma, t > 0,
\]

\[
\check{u}(0) = \check{u}_0 \quad \text{in } \Omega \setminus \Sigma, \quad \check{h}(0) = \check{h}_0 \quad \text{on } \Sigma,
\]

with \(\check{z}_0 = (\check{u}_0, \check{h}_0) = (\phi(\check{u}_0, \check{h}_0), 0)\) and \(E(\bar{u}) := \nabla \bar{u} + [\nabla \bar{u}]^T\). Writing this problem abstractly as \(L_{\omega, \bar{z}} = N(\bar{z} + \check{z})\), by the implicit function theorem we obtain a unique solution \(\bar{z} = \bar{z}(\check{z})\) in the function space \(E(a)\) for each \(a < t_\ast\). Then \(\check{z}\) is determined by the problem

\[
\begin{align*}
\rho \partial_t \check{u} - \mu \Delta \check{u} + \nabla \check{\pi} &= \rho \omega \check{u} & \text{in } \Omega \setminus \Sigma, t > 0, \\
\text{div } \check{u} &= 0 & \text{in } \Omega \setminus \Sigma, t > 0, \\
-P_{\Sigma} [\mu (\nabla \check{u} + [\nabla \check{u}]^T)]\nu_{\Sigma} &= 0 & \text{on } \Sigma, t > 0, \\
-([\mu (\nabla \check{u} + [\nabla \check{u}]^T)]\nu_{\Sigma}) + [\check{\pi}] - \sigma A_{\Sigma} \check{h} &= G_\gamma(\check{h}) & \text{on } \Sigma, t > 0, \\
[\check{u}] &= 0 & \text{on } \Sigma, t > 0, \\
\check{u} &= 0 & \text{on } \partial \Omega, t > 0,
\end{align*}
\]

\[
\check{u}(0) = \check{u}_0 \quad \text{in } \Omega \setminus \Sigma, \quad \check{h}(0) = \check{h}_0 \quad \text{on } \Sigma.
\]

The last equation can be rewritten abstractly in \(X_0\) employing the operator \(A\) introduced above as

\[
\dot{\check{z}} + A\check{z} = R(\check{z}), \quad t > 0, \quad \check{z}(0) = \check{z}_0,
\]

where

\[
R(\check{z}) = (\omega(I - T_1)\check{u}(\check{z}) - T_2 G_\gamma(\check{z}), \omega \check{h}(\check{z})).
\]

Note that \(\check{z}\) is a causal functional of \(\check{z}\).

3. Problem (6.9) is of the form studied in [22], where the generalized principle of linearized stability is proved for abstract parabolic quasilinear problems of the form (6.9). The only difference is that here a part of \(R\) is nonlocal, but causal in time. Therefore we only comment on the required modifications in the proof of Theorem 2.1 in [22]. For this purpose we decompose

\[
R(\check{z}) = R_{nloc}(\check{z}) + R_{loc}(\check{z}) := (\omega(I - T_1)\check{u}(\check{z}), \omega \check{h}(\check{z})) + (-T_2 G_\gamma(\check{z}), 0).
\]

Observe that by construction, if \(z = \check{z}\) is an equilibrium then \(z = \check{z}\), hence \(\check{z} = 0\). Therefore the equilibria are determined by the equation \(Az_\ast = R_{loc}(z_\ast)\). Further, we have an estimate of the form

\[
|R_{nloc}(\check{z})|_{E(t_0)} \leq \varepsilon|\check{u}|_{E(t_0)},
\]
provided $\tilde{z}_0$ is small in the norm of $X_\gamma$.

Let $P^c$ denote the projection in $X_0$ onto the kernel $N(A)$ along the range $R(A)$ of $A$ and let $P^s = I - P^c$ the complementary projection onto $R(A)$. As in the proof of Theorem 2.1 of [22] we parameterize the set of equilibria $E$ near 0 over $N(A)$ via a $C^1$-map $x \mapsto x + \Psi(x)$ such that $\Psi(0) = \Psi'(0) = 0$. By $x := P^c \tilde{z}$ and $y := P^s \tilde{z} - \Psi(P^c \tilde{z})$ we introduce the normal form of the problem which reads as follows.

$$
\begin{align*}
\dot{x} &= T(x, y), \quad x(0) = x_0, \tag{6.10} \\
\dot{y} + A^* y &= S(x, y), \quad y(0) = y_0,
\end{align*}
$$

where

$$
T(x, y) = P^c [R(x + \Psi(x) + y) - R(x + \Psi(x))],
$$

$$
S(x, y) = P^s [R(x + \Psi(x) + y) - R(x + \Psi(x))] - \Psi'(x) T(x, y),
$$

since $P^c R(x + \Psi(x)) = 0$ and $P^s R(x + \Psi(x)) = A^* \Psi(x) = A^*[x + \Psi(x)]$. The first component of $P^c \tilde{z}$ equals zero, since the eigenfunctions of $A$ for eigenvalue 0 have vanishing velocity part. This implies $|\tilde{u}|_{E, t} \leq |y|_{E(t)}$, hence the nonlocal part is estimated as

$$
|R_{\text{loc}}(\tilde{z})|_{E(t)} \leq \varepsilon |\tilde{u}|_{E, t} \leq \varepsilon |y|_{E(t)}
$$

for any $t_0 > 0$. The local parts of $T$ and $R$ can be estimated by $\varepsilon |y|_{E(t)}$, as in the proof of Theorem 2.1 in [22]. Taking these observations into account, we may proceed as in [22] to prove global existence of $\tilde{z}$, its stability and convergence to another equilibrium, provided $\tilde{z}_0$ is small in $X_\gamma$.

\[\square\]

7. Global Existence and Convergence

Again we assume for simplicity that the phases are connected. There are basically two obstructions against global existence:

- **regularity:** the norms of either $u(t)$ or $\Gamma(t)$ become unbounded;
- **geometry:** the topology of the interface changes, or the interface touches the boundary of $\Omega$.

Note that the *phase volumes* are preserved by the semiflow.

We say that a solution $(u, \Gamma)$ satisfies a **uniform ball condition**, if there is a radius $r > 0$ such that $\Gamma([0, t_\ast)) \subset \mathcal{M}H^2(\Omega, r)$. Note that this condition bounds the curvature of $\Gamma(t)$, and prevents it to touch the outer boundary $\partial G$, or to undergo topological changes.

Combining the above results, we obtain the following result on the asymptotic behavior of solutions.

**Theorem 7.1.** Let $p > n + 2$. Suppose that $(u, \Gamma)$ is a solution of the two-phase Navier-Stokes problem with surface tension on the maximal time interval $[0, t_\ast)$.

Assume the following on $[0, t_\ast)$:

(i) $|u(t)|_{W^{2-2/p, p}_p} + |\Gamma(t)|_{W^{3-2/p, p}_p} \leq M < \infty$;

(ii) $(u, \Gamma)$ satisfies a uniform ball condition.

Then $t_\ast = \infty$, i.e. the solution exists globally, and it converges in $PM$ to an equilibrium $(0, \Gamma_\infty) \in E$. Conversely, if $(u, \Gamma)$ is a global solution in $PM$ which converges to an equilibrium $(u_\infty, \Gamma_\infty) \in E$ in $PM$, then (i) and (ii) are valid.
Proof. Assume that (i) and (ii) are valid. Then $\Gamma([0,t_\ast)) \subset W^{3-2/p}_{p}(\Omega,r)$ is bounded, hence relatively compact in $W^{3-2/p-\varepsilon}_{p}(\Omega,r)$. Thus we may cover $\Gamma([0,t_\ast))$ by finitely many balls with centers $\Sigma_k$ such that $\text{dist}_{W^{3-2/p-\varepsilon}_{p}}(\Gamma(t),\Sigma_j) \leq \delta$ for some $j = j(t), \ t \in [0,t_\ast)$. Let $J_k = \{ t \in [0,t_\ast) : j(t) = k \}$; using for each $k$ a Hanzawa-transformation $\Theta_k$, we see that the pull backs $\{ u(t,) \circ \Theta_k : t \in J_k \}$ are bounded in $W^{3-2/p}_{p}(\Omega\setminus\Sigma_k)$, hence relatively compact in $W^{3-2/p-\varepsilon}_{p}(\Omega\setminus\Sigma_k)$. Employing now Corollary 4.4 with $\mu = 1-\varepsilon$ we obtain solutions $(u^1,\Gamma^1)$ with initial configurations $(u(t),\Gamma(t))$ in the phase manifold on a common time interval say $(0,a]$, and by uniqueness we have $(u^1(a),\Gamma^1(a)) = (u(t+a),\Gamma(t+a))$. Continuous dependence implies then relative compactness of $\{ (u(t),\Gamma(t)) : 0 \leq t < t_\ast \}$ in $PM$, in particular $t_\ast = \infty$ and the orbit $(u,\Gamma)(\mathbb{R}_+) \subset PM$ is relatively compact. The energy is a strict Ljapunov functional, hence the limit set $\omega(u,\Gamma)$ of a solution is contained in the set $E$ of equilibria. By compactness $\omega(u,\Gamma) \subset PM$ is non-empty, hence the solution comes close to $E$. Finally, we apply the convergence result Theorem 6.3 to complete the sufficiency part of the proof. Necessity follows by a compactness argument. \hfill $\square$

8. Appendix: Transmission Problems

In this section we provide some results, concerning the existence and uniqueness of solutions to the transmission problem

$$\lambda q - \Delta q = f, \quad x \in \Omega \setminus \Gamma$$

$$[\partial_{\nu}q] = g, \quad x \in \Gamma,$$

$$[\partial_{\nu\nu}q] = h_1, \quad x \in \Gamma,$$

$$\delta \partial_{\nu\nu\nu}q + (1-\delta)q = h_{2,\delta}, \quad x \in \partial\Omega, \delta \in \{0,1\},$$

where $\lambda \geq 0$,

$$\rho(x) := \rho_1 \chi_{\Omega_1}(x) + \rho_2 \chi_{\Omega_2}(x), \quad x \in \Omega \setminus \Gamma,$$

and $\rho_j > 0$. To be precise, we will study (8.1) in different functional analytic settings. We begin by stating the result for the `classical' case, i.e. if the basic space is given by $L_p(\Omega)$.

**Theorem 8.1.** Let $\Omega \subset \mathbb{R}^n$ open, $1 < p < \infty$, $f \in L_p(\Omega)$, $g \in W^{2-1/p}_{p}(\Gamma)$, $h_1 \in W^{1-1/p}_{p}(\Gamma)$ and $h_{2,\delta} \in W^{2-\delta-1/p}_{p}(\partial\Omega)$, $\delta \in \{0,1\}$ be given. Then, for each $\lambda > 0$, there exists a unique solution $q \in H^2_p(\Omega\setminus\Gamma)$ of (8.1) and a constant $C_1 > 0$ such that

$$|q|_{H^2_p(\Omega\setminus\Gamma)} \leq C_1 \left( |f|_{L_p(\Omega)} + |g|_{W^{2-1/p}_{p}(\Gamma)} + |h_1|_{W^{1-1/p}_{p}(\Gamma)} + |h_{2,\delta}|_{W^{2-\delta-1/p}_{p}(\partial\Omega)} \right).$$

If in addition $J = [0,a]$, $f = f(t,x)$, $f \in H^1_p(J;L_p(\Omega))$ and $g = h_1 = h_{2,\delta} = 0$, then for each $\lambda > 0$, there exists a unique solution $q \in H^2_p(J;H^1_p(\Omega\setminus\Gamma))$, and the estimate

$$||q||_{H^2_p(J;H^2_p(\Omega\setminus\Gamma))} \leq C_2 ||f||_{H^1_p(J;L_p(\Omega))}$$

holds with some constant $C_2 > 0$.

**Proof.** The first assertion basically follows from [12], since the Lopatinskii-Shapiro condition is satisfied at $\Gamma$ and $\partial\Omega$. The second assertion follows from the first one by differentiating (8.1) w.r.t. $t$ and by employing the uniqueness of the solution of (8.1). \hfill $\square$
We will also need a result for the case \( \lambda = 0 \). To this end, let \( \Omega \subset \mathbb{R}^n \) be a bounded domain, \( g = h_1 = h_2, \delta = 0 \) and \( f \in L_p(\Omega) \). Define \( A_\delta \) by \( A_\delta q = -\Delta q \), with domain

\[
D(A_\delta) = \{ q \in H^2_p(\Omega \setminus \Gamma) : [\rho q] = [\partial_x q] = 0 \text{ on } \Gamma, \\
\delta \partial_{\nu_\delta} q + (1 - \delta)q = 0, \text{ on } \partial \Omega \}, \; \delta \in \{0, 1\}.
\]

Since

\[
D(A_\delta) \hookrightarrow L_p(\Omega),
\]
the resolvent of \( A_\delta \) is compact and therefore the spectral set \( \sigma(A_\delta) \) consists solely of a countably infinite sequence of isolated eigenvalues. In case \( \delta = 1 \) it can be readily checked that \( 0 \) is a simple eigenvalue of \( A_1 \), hence \( L_p(\Omega) = N(A_1) \oplus R(A_1) \).

The kernel \( N(A_1) \) of \( A_1 \) is given by

\[
N(A_1) = \mathbb{K}\mathbb{I}_p, \quad \text{where}
\]

\[
\mathbb{I}_p(x) := \chi_{\Omega_1}(x) + \rho_1 \rho_2 \chi_{\Omega_2}(x), \quad x \in \Omega \setminus \Gamma.
\]

and \( R(A_1) = \{ f \in L_p(\Omega) : (f|\mathbb{I}_p) = 0 \} \). Therefore \( (8.1) \) has a unique solution \( q \in H^2_p(\Omega) \setminus \Gamma) \oplus \mathbb{K}\mathbb{I}_p \), provided \( (f|\mathbb{I}_p) = 0 \). In case of Dirichlet boundary conditions, i.e. \( \delta = 0 \), it holds that \( N(A_0) = \{0\} \), hence or each \( f \in L_p(\Omega) \), the system \( (8.1) \) admits a unique solution \( q \in H^1_p(\Omega \setminus \Gamma) \).

**Theorem 8.2.** Let \( \Omega \subset \mathbb{R}^n \) a bounded domain, \( 1 < p < \infty, f \in L_p(\Omega), \; g = h_1 = h_2 = 0 \) and \( \lambda = 0 \). Then the following assertions hold

1. If \( \delta = 0 \), then there exists a unique solution \( q \in H^2_p(\Omega \setminus \Gamma) \) of \( (8.1) \).
2. If \( \delta = 1 \) and \( (f|\mathbb{I}_p) = 0 \), then there exists a unique solution \( q \in H^2_p(\Omega) \setminus \Gamma) \oplus \mathbb{K}\mathbb{I}_p \).

If in addition \( J = [0,a], \; f = f(t,x) \) and \( f \in H^1_p(J; L_p(\Omega)) \) s.t. \( f(t,\cdot) \in R(A_\delta) \) for a.e. \( t \in J \), then \( q \in H^1_p(J; H^2_p(\Omega) \setminus \Gamma) \oplus N(A_\delta) \).

8.1. **A weak transmission problem.** Here we study the (weak) transmission problem

\[
(\nabla q | \nabla \phi)_{L_2(\Omega)} = (f|\nabla \phi)_{L_2(\Omega)}, \quad \phi \in H^1_p(\Omega), \\
[\rho q] = g, \quad x \in \Gamma,
\]

where \( \Omega \subset \mathbb{R}^n \) is open and bounded with \( \partial \Omega \in C^2 \). We want to show that this problem admits a unique solution \( q \in H^1_p(\Omega \setminus \Gamma) \), that satisfies the estimate

\[
|\nabla q|_{L_p(\Omega)} \leq C \left( |f|_{L_p(\partial \Omega; \mathbb{R}^n)} + |g|_{W^{2-1/p}_p(\Gamma)} \right),
\]

provided \( f \in L_p(\Omega; \mathbb{R}^n) \) and \( g \in W^{2-1/p}_p(\Gamma) \). We will first treat the case \( f = 0 \), \( g \in W^{2-1/p}_p(\Gamma) \) and consider the problem

\[
\lambda(q|\phi) + (\nabla q | \nabla \phi)_{L_2(\Omega)} = 0, \quad \phi \in H^1_p(\Omega), \\
[\rho q] = g, \quad x \in \Gamma.
\]

with \( \lambda > 0 \). Theorem 8.2 then yields a strong unique solution \( q \in H^2_p(\Omega \setminus \Gamma) \) of \( (8.1) \) with \( f = h_1 = h_2 = 0 \) which is also the unique solution of \( (8.2) \). This follows from integration by parts. Our aim is to derive an estimate which is of the form

\[
|g|_{H^1_p(\Omega \setminus \Gamma)} \leq C |g|_{W^{2-1/p}_p(\Gamma)},
\]
which will be done by a localization argument. For this purpose we consider first the following auxiliary transmission problem

\[ \lambda q - \Delta q = f, \quad x \in \mathbb{R}^n, \]
\[ [\rho q] = g, \quad x \in \mathbb{R}^{n-1}, \]
\[ [\partial_x q] = h, \quad x \in \mathbb{R}^{n-1}, \]

with data \( f \in L_p(\mathbb{R}^n), \) \( g \in W^{2-1/p}_p(\mathbb{R}^{n-1}) \) and \( h \in W^{1-1/p}_p(\mathbb{R}^n) \), which will play an important role in the forthcoming localization procedure. Solve the full space problem

\[ \lambda q - \Delta q = f, \quad x \in \mathbb{R}^n, \]

and obtain a unique solution \( q_1 = (\lambda - \Delta)^{-1} f \in H^2_p(\mathbb{R}^n) \), provided \( \text{Re}\lambda > 0 \). In the sequel we will always assume that \( \lambda \) is real and \( \lambda \geq 1 \). In particular, it follows that

\[ \lambda^{1/2}|q_1|_{L_p(\mathbb{R}^n)} + |\nabla q_1|_{L_p(\mathbb{R}^n)} \leq C|f|_{H^{-1}_p(\mathbb{R}^n)}, \] (8.4)

with some constant \( C > 0 \) being independent of \( \lambda \geq 1 \), since

\[ \lambda^{1/2}|(\lambda - \Delta)^{-1} f|_{L_p(\mathbb{R}^n)} \leq C\lambda^{1/2}|||I - \Delta||^1/2(\lambda - \Delta)^{-1} |||B(L_p,L_p)|| f|_{H^{-1}_p(\mathbb{R}^n)} \]
\[ \leq C|||I - \Delta||^{1/2}(\lambda - \Delta)^{-1/2} |||B(L_p,L_p)|| f|_{H^{-1}_p(\mathbb{R}^n)} \]
\[ \leq C|f|_{H^{-1}_p(\mathbb{R}^n)}, \]

and

\[ |\nabla(\lambda - \Delta)^{-1} f|_{L_p(\mathbb{R}^n)} \leq C||(I - \Delta)(\lambda - \Delta)^{-1} |||B(L_p,L_p)|| f|_{H^{-1}_p(\mathbb{R}^n)} \]
\[ \leq C|f|_{H^{-1}_p(\mathbb{R}^n)}, \]

since the norm

\[ |||I - \Delta||^\alpha(\lambda - \Delta)^{-\alpha}|||B(L_p,L_p)||, \quad \alpha \in \{1/2, 1\}, \]

is independent of \( \lambda \geq 1 \), which follows e.g. from functional calculus. The shifted function \( q_2 = q - q_1 \) should now solve the reduced problem

\[ \lambda q_2 - \Delta q_2 = 0, \quad x \in \mathbb{R}^n, \]
\[ [\rho q_2] = \tilde{g}, \quad x \in \mathbb{R}^{n-1}, \] (8.5)
\[ [\partial_x q_2] = h, \quad x \in \mathbb{R}^{n-1}, \]

with a modified function \( \tilde{g} \in W^{2-1/p}_p(\mathbb{R}^{n-1}). \) Let \( x = (x', y) \in \mathbb{R}^n \times \mathbb{R} \) and define \( L := (\lambda - \Delta_n)^{1/2}, \) where \( \Delta_n \) denotes the Laplacian with respect to the first \( n - 1 \) variables \( x' \) and with domain \( D(L) = H^1_p(\mathbb{R}^{n-1}). \) Let furthermore

\[ \rho(x', y) = \rho_2 \chi_{\{y > 0\}}(x', y) + \rho_1 \chi_{\{y < 0\}}(x', y), \quad (x', y) \in \mathbb{R}^{n-1} \times \mathbb{R}. \]

We make the following ansatz to find a solution of (8.5)

\[ q_2(y) := \begin{cases} e^{-Ly}a_+, & y > 0, \\ e^{Ly}a_-, & y < 0, \end{cases} \] (8.6)

where \( a_-, a_+ \) have to be determined. The first transmission condition in (8.5) yields \( \rho_2 a_+ - \rho_1 a_- = \tilde{g} \), whereas the second condition implies \( -L(a_+ + a_-) = h \), hence
a_+ + a_- = -L^{-1}h. Observe that \( \tilde{g}, L^{-1}h \in W^{2-1/p}_{p}(\mathbb{R}^{n-1}) \). Therefore we may solve this linear system of equations to the result

\[
a_- = -\frac{1}{\rho_1 + \rho_2} (\tilde{g} + \rho_2 L^{-1}h), \quad a_+ = \frac{1}{\rho_1 + \rho_2} (\tilde{g} + \rho_2 L^{-1}h) - L^{-1}h. \tag{8.7}
\]

In other words, the solution of (8.5) (hence of (8.3)) is uniquely determined and \( a_- , a_+ \in W^{2-1/p}_{p}(\mathbb{R}^{n-1}) \). Since \( |L^{\pm} \cdot |_{L^p(\mathbb{R}^{n-1} \times \mathbb{R}_+)} \) is an equivalent norm in \( W^{1-1/p}_{p}(\mathbb{R}^{n-1}) \) and the corresponding constants are independent of \( \lambda \geq 1 \), we obtain first

\[
\lambda^{1/2} |q_2|_{L^p(\mathbb{R}^n)} = \lambda^{1/2} |L^{-1}Lq_2|_{L^p(\mathbb{R}^n)} \leq C \left( |a_+|_{W^{1-1/p}_{p}(\mathbb{R}^{n-1})} + |a_-|_{W^{1-1/p}_{p}(\mathbb{R}^{n-1})} \right).
\]

Concerning \( \nabla q_2 \) in \( L^p(\mathbb{R}^n) \), we estimate as follows

\[
|\nabla_x q_2|_{L^p(\mathbb{R}^n)} \leq C |L_0 q_2|_{L^p(\mathbb{R}^n)} = C |L_0 L^{-1}Lq_2|_{L^p(\mathbb{R}^n)} \\
\leq C \|L_0 L^{-1}\|_{B(L^p, L^p)} |Lq_2|_{L^p(\mathbb{R}^n)} \\
\leq C \left( |a_+|_{W^{1-1/p}_{p}(\mathbb{R}^{n-1})} + |a_-|_{W^{1-1/p}_{p}(\mathbb{R}^{n-1})} \right),
\]

with \( L_0 := (I - \Delta_x)^{1/2} \). Here the norm \( \|L_0 L^{-1}\|_{B(L^p, L^p)} \) does not depend on \( \lambda \geq 1 \), which is a consequence of the functional calculus. The estimate for \( \partial_y q_2 \) in \( L^p(\mathbb{R}^n) \) is even simpler, since

\[
|\partial_y q_2|_{L^p(\mathbb{R}^n)} = |Lq_2|_{L^p(\mathbb{R}^n)} \leq C \left( |a_+|_{W^{1-1/p}_{p}(\mathbb{R}^{n-1})} + |a_-|_{W^{1-1/p}_{p}(\mathbb{R}^{n-1})} \right).
\]

This yields the estimate

\[
\lambda^{1/2} |q_2|_{L^p(\mathbb{R}^n)} + |\nabla q_2|_{L^p(\mathbb{R}^n)} \leq C \left( |\tilde{g}|_{W^{1-1/p}_{p}(\mathbb{R}^{n-1})} + |L^{-1}h|_{W^{1-1/p}_{p}(\mathbb{R}^{n-1})} \right).
\]

For each fixed \( \lambda \geq 1 \) the operator \( L^{-1} \) is bounded and linear from \( W^{1-1/p}_{p}(\mathbb{R}^{n-1}) \) to \( W^{1-1/p}_{p}(\mathbb{R}^{n-1}) \), where \( W^{1-1/p}_{p}(\mathbb{R}^{n-1}) \) is the topological dual space of \( W^{p'}_{1}(\mathbb{R}^{n-1}) \), and \( 1/p + 1/p' = 1 \). We want to show that the bound of \( L^{-1} \) is independent of \( \lambda \geq 1 \). This can be seen as follows. We have

\[
|L^{-1}h|_{W^{1}_{p}(\mathbb{R}^{n-1})} \leq C |L_0 L^{-1}h|_{L^p(\mathbb{R}^{n-1})} \leq C \|L_0 L^{-1}\|_{B(L^p, L^p)} |h|_{L^p(\mathbb{R}^{n-1})},
\]

which holds for all \( h \in L^p(\mathbb{R}^{n-1}) \), since \( |L_0 \cdot |_{L^p(\mathbb{R}^{n-1})} \) is an equivalent norm in \( W^{1}_{p}(\mathbb{R}^{n-1}) \). On the other hand we have

\[
|L^{-1}h|_{L^p(\mathbb{R}^{n-1})} = |L_0 L_0^{-1} L^{-1}h|_{L^p(\mathbb{R}^{n-1})} = |L_0 L^{-1}L_0^{-1}h|_{L^p(\mathbb{R}^{n-1})} \\
\leq \|L_0 L^{-1}\|_{B(L^p, L^p)} |L_0^{-1}h|_{L^p(\mathbb{R}^{n-1})} \\
\leq C \|L_0 L^{-1}\|_{B(L^p, L^p)} |h|_{W^{1}_{p}(\mathbb{R}^{n-1})},
\]

for all \( h \in W^{1}_{p}(\mathbb{R}^{n-1}) \), since \( |L_0^{-1} \cdot |_{L^p(\mathbb{R}^{n-1})} \) is an equivalent norm in \( W^{p'}_{1}(\mathbb{R}^{n-1}) \) and since \( L^{-1} \) and \( L_0^{-1} \) are commuting operators. Finally we apply the real interpolation method to obtain

\[
|L^{-1}h|_{W^{1-1/p}_{p}(\mathbb{R}^{n-1})} \leq C |h|_{W^{1-1/p}_{p}(\mathbb{R}^{n-1})},
\]

for all \( h \in W^{1-1/p}_{p}(\mathbb{R}^{n-1}) \), where the constant \( C > 0 \) is independent of \( \lambda \geq 1 \). In summary we derived the a priori estimate

\[
\lambda^{1/2} |q_2|_{L^p(\mathbb{R}^n)} + |\nabla q_2|_{L^p(\mathbb{R}^n)} \leq C \left( |\tilde{g}|_{W^{1-1/p}_{p}(\mathbb{R}^{n-1})} + |h|_{W^{1-1/p}_{p}(\mathbb{R}^{n-1})} \right),
\]
for the solution of (8.5), hence
\[
\lambda^{1/2} |q|_{L_p(\mathbb{R}^n)} + |\nabla q|_{L_p(\mathbb{R}^n)}
\leq C \left( |f|_{H_p^{-1}(\mathbb{R}^n)} + |g|_{W_p^{1-1/p}(\mathbb{R}^{n-1})} + |h|_{W_p^{1-1/p}(\mathbb{R}^{n-1})} \right) \tag{8.8}
\]
for the solution of (8.3), since
\[
|\bar{g}|_{W_p^{1-1/p}(\mathbb{R}^{n-1})} \leq |g|_{W_p^{1-1/p}(\mathbb{R}^{n-1})} + ||\rho q||_{W_p^{1-1/p}(\mathbb{R}^{n-1})}
\leq |g|_{W_p^{1-1/p}(\mathbb{R}^{n-1})} + C|f|_{H_p^{-1}(\mathbb{R}^n)}.
\]
by (8.1). Consider now a bounded domain \( \Omega \subset \mathbb{R}^n \) with \( \partial \Omega \in C^2 \) and let \( \Gamma \subset \Omega \) be a hypersurface such that \( \Gamma \) divides the set \( \Omega \) into two disjoint regions \( \Omega_1, \Omega_2 \), where \( \partial \Omega_1 = \Gamma \) and \( \partial \Omega_2 = \partial \Omega \cup \Gamma \). Since \( \Omega \) is compact, we may cover it by a union of finitely many open sets \( U_k, k = 0, \ldots, N \) which are subject to the following conditions
- \( \partial \Omega \subset U_0 \) and \( U_0 \cap \Gamma = \emptyset \);
- \( U_1 \subset \Omega_1 \) and \( U_1 \cap \Gamma = \emptyset \);
- \( U_k \cap \Gamma \neq \emptyset, U_k \cap \partial \Omega = \emptyset \) for \( k = 2, \ldots, N \) and
\[
\bigcup_{k=2}^N U_k \supset \Gamma.
\]
For \( k \geq 2 \), the sets \( U_k \) may be balls with a fixed but arbitrarily small radius \( r > 0 \).
Let \( \{ \varphi_k \}_{k=0}^N \) be a partition of unity, such that \( \text{supp} \varphi_k \subset U_k \) and \( 0 \leq \varphi_k(x) \leq 1 \) for all \( x \in \Omega \). Consider the transmission problem
\[
\begin{align*}
\lambda q - \Delta q &= 0, \quad x \in \Omega \setminus \Gamma \\
[\rho q] &= g, \quad x \in \Gamma, \\
[\partial_n q] &= 0, \quad x \in \Gamma, \\
\partial_n q &= 0, \quad x \in \partial \Omega,
\end{align*}
\tag{8.9}
\]
where \( g \in W_p^{2-1/p}(\Gamma) \). Set \( q_k = \varphi_k q \) and \( g_k = \varphi_k g \). By Theorem (8.1) there exists a unique solution \( q \in H_p^2(\Omega \setminus \Gamma) \) of (8.9), if e.g. \( \lambda \geq 1 \). Multiplying (8.9) by \( \varphi_0 \) yields
\[
\begin{align*}
\lambda q_0 - \Delta q_0 &= -2(\nabla q_0 \nabla \varphi_0) - q \Delta \varphi_0, \quad x \in \Omega, \\
\partial_n q_0 &= q \partial_n \varphi_0, \quad x \in \partial \Omega,
\end{align*}
\tag{8.10}
\]
which is an elliptic boundary value problem in \( \Omega \). Denote by \( (F_0, G_0) \) the right hand side of (8.10). By [29] Theorem 3.3.4, there exists a common bounded extension operator \( E \) from \( L_p(\Omega) \) resp. \( H_p^{-1}(\Omega) \) to \( L_p(\mathbb{R}^n) \) resp. \( H_p^{-1}(\mathbb{R}^n) \). Solve the equation
\[
\lambda q_0^1 - \Delta q_0^1 = EF_0, \quad x \in \mathbb{R}^n.
\]
The solution is given by \( q_0^1 = (\lambda - \Delta)^{-1} EF_0 \) and we have the estimate
\[
\lambda^{1/2} |q_0|_{L_p(\mathbb{R}^n)} + |\nabla q_0|_{L_p(\mathbb{R}^n)} \leq C|EF_0|_{H_p^{-1}(\mathbb{R}^n)} \leq C|F_0|_{H_p^{-1}(\Omega)} \leq C|q|_{L_p(\Omega)},
\]
as we have already shown. Note that since \( F_0 \in L_p(\Omega) \), it holds that
\[
|q_0^1|_{H_p^2(\mathbb{R}^n)} = |(\lambda - \Delta)^{-1} EF_0|_{H_p^2(\mathbb{R}^n)} \leq C|F_0|_{L_p(\Omega)} \leq C|q|_{H_p^2(\Omega)},
\]
and \( C > 0 \) does not depend on \( \lambda \geq 1 \). In particular, the real interpolation method yields
\[
|q_0^1|_{W_p^{s+\epsilon}(\mathbb{R}^n)} \leq C|q|_{W_p^s(\Omega)}, \quad s \in [0, 1].
\]
The shifted function $q_0^2 = q_0 - q_0^1$ solves the problem
\[
\lambda q_0^2 - \Delta q_0^2 = 0, \quad x \in \Omega,
\]
\[
\partial_{n\nu} q_0^2 = G_0^2, \quad x \in \partial \Omega,
\]
with some modified function $G_0^2 \in W_p^{1-1/p}(\partial \Omega)$. By [1, Theorem 9.2], there exists a bounded solution operator $S_0^2 : W_p^{-1/p}(\partial \Omega) \to H_p^1(\Omega)$ such that $q_0^2 = S_0^2 G_0^2$ and there exists a constant $C > 0$ independent of $\lambda \geq 1$ such that
\[
\lambda^{1/2} |q_0^2|_{L_p(\Omega)} + |\nabla q_0^2|_{L_p(\Omega)} \leq C |G_0^2|_{W_p^{-1/p}(\partial \Omega)}.
\]
This yields
\[
\lambda^{1/2} |q_0|_{L_p(\Omega)} + |\nabla q_0|_{L_p(\Omega)} \leq C \left( |\nabla q_0|_{H_p^{-1}(\Omega)} + |q_0|_{H_p^{-1}(\Omega)} \right.
\]
\[
\left. + |q_0|^s_{W_p^{-1/p}(\partial \Omega)} + |\partial_{n\nu} q_0|_{W_p^{-1/p}(\partial \Omega)}. \right)
\]
Since $\varphi_0$ is smooth and compactly supported and since $\nu \in C^1(\partial \Omega)$, we have
\[
\lambda^{1/2} |q_0|_{L_p(\Omega)} + |\nabla q_0|_{L_p(\Omega)} \leq C |q_0|_{W_p^2(\Omega)},
\]
for some $s \in (1/p, 1)$, since
\[
|q_0|^s_{W_p^{-1/p}(\partial \Omega)} \leq C |q_0|_{L_p(\partial \Omega)} \leq C |q_0|_{W_p^2(\Omega)}, \quad s \in (1/p, 1),
\] and
\[
|\partial_{n\nu} q_0|^s_{W_p^{-1/p}(\partial \Omega)} \leq C |\partial_{n\nu} q_0|_{L_p(\partial \Omega)} \leq C |q_0|^s_{W_p^2(\Omega)} \leq C |q_0|_{W_p^2(\Omega)}.
\]
In a next step we multiply (8.9) by $\varphi_1$ to obtain the full space problem
\[
\lambda q_1 - \Delta q_1 = -2(\nabla q \nabla \varphi_1) - q \Delta \varphi_1, \quad x \in \mathbb{R}^n.
\]
This problem admits a unique solution $q_1 = (\lambda - \Delta)^{-1} F_1$, provided $\lambda \geq 1$, where $S_1 = (\lambda - \Delta)^{-1} : H_p^{-1}(\mathbb{R}^n) \to H_p^1(\mathbb{R}^n)$ is bounded and $F_1$ denotes the right hand side of (8.13). As before we obtain the estimate
\[
\lambda^{1/2} |q_1|_{L_p(\mathbb{R}^n)} + |\nabla q_1|_{L_p(\mathbb{R}^n)} \leq C |q_1|_{L_p(\Omega)},
\]
with $C > 0$ being independent of $\lambda \geq 1$.

We turn now to the charts $U_k$, $k = 2, \ldots, N$. Multiplying (8.9) by $\varphi_k$, $k = 2, \ldots, N$, we obtain the pure transmission problem
\[
\lambda q_k - \Delta q_k = -2(\nabla q \nabla \varphi_k) - q \Delta \varphi_k, \quad x \in \mathbb{R}^n \setminus \Gamma,
\]
\[
[\partial_{n\nu} q_k] = g_k, \quad x \in \Gamma,
\]
\[
[\partial_{n\nu} \varphi_k] = [q] \partial_{n\nu} \varphi_k, \quad x \in \Gamma.
\]
Let $x_0 \in \Gamma$. Then there exists $k \in \{2, \ldots, N\}$ such that $x_0 \in U_k$. After a translation and a rotation of coordinates we may assume that $x_0 = 0$ and that the normal $\nu(x_0)$ at $x_0$ which points from $\Omega_1$ to $\Omega_2$ is given by $\nu(x_0) = [0, \ldots, 0, -1]^T$. Consider a graph $\eta \in C^2(\mathbb{R}^{n-1})$ with compact support such that
\[
\{(x', x_n) \in U_k \subset \mathbb{R}^{n-1} \times \mathbb{R} : x_n = \eta(x')\} = \Gamma \cap U_k.
\]
Note that, since $\nabla x' \eta(0) = 0$, we may choose $|\nabla x' \eta|_\infty$ as small as we wish, by decreasing the size of the chart $U_k$. Let $q(x', x_n) = \nu(x', x_n - \eta(x'))$, where $(x', x_n) \in \{(x', x_n) \in U_k \subset \mathbb{R}^{n-1} \times \mathbb{R} : x_n = \eta(x')\} = \Gamma \cap U_k$. 


declamation
$U_k$. We define a new coordinate by $y = x_n - \eta(x'), (x', x_n) \in U_k$. Then we obtain
\[
\Delta q(x', x_n) = \Delta y v(x', y) - 2\partial_{y^i} (\nabla_{x'} v(x', y) \nabla_{x'} \eta(x')) + \partial_{y^j}^2 v(x', y) |\nabla_{x'} \eta|^2 - \partial_{y^j} v(x', y) \Delta_{x'} \eta(x')
\]
and
\[
[\partial_y q] = -\frac{1}{\sqrt{1 + |\nabla_{x'} \eta|^2}} [\partial_{y^i} v] + \frac{1}{\sqrt{1 + |\nabla_{x'} \eta|^2}} ([\nabla_{x'} v]|\nabla_{x'} \eta),
\]
since the normal at $x \in U_k \cap \Gamma$ is given by
\[
\nu(x', \eta(x')) = \frac{1}{\sqrt{1 + |\nabla_{x'} \eta|^2}} ([\nabla_{x'} v]|\nabla_{x'} \eta)^T, -1]^T.
\]
Let $(\Theta u)(x', y) := q(x', y + \eta(x')) = v(x', y)$ with inverse $(\Theta^{-1} v)(x', x_{n+1}) = v(x', x_{n+1} - \eta(x')) = q(x', x_{n+1})$. Applying the $C^2$-diffeomorphism $\Theta$ to (8.15) and considering the terms on the right hand side of (8.16) which depend on $u$ as given functions $(f_k, g_k, h_k)$ yields the problem
\[
\lambda v_k - \Delta v_k = F(f_k, v_k, \varphi_k, \eta), \quad (x', y) \in \mathbb{R}^n,
\]
\[
[\rho v_k] = G(g_k), \quad x' \in \mathbb{R}^{n-1}, \quad y = 0,
\]
\[
[\partial_{y^i} v_k] = H(v_k, \varphi_k, \eta), \quad x' \in \mathbb{R}^{n-1}, \quad y = 0.
\]
which is of the form (8.3). Here
\[
F(f_k, v_k, \varphi_k, \eta) := -2\partial_{y^i} (\nabla_{x'} v_k |\nabla_{x'} \eta) + \partial_{y^i}^2 v_k |\nabla_{x'} \eta|^2 - \partial_{y^i} v_k \Delta_{x'} \eta,
\]
\[
G(g_k) := \Theta g_k
\]
and
\[
H(h_k, v_k, \varphi_k, \eta) := \frac{1}{1 + |\nabla_{x'} \eta|^2} ([\nabla_{x'} v_k]|\nabla_{x'} \eta)
\]
We want to apply (8.8) to (8.16) and estimate as follows.
\[
|\partial_{y^i} (\nabla_{x'} v_k |\nabla_{x'} \eta)|_{W^{-1}_{p-1}([\mathbb{R}^n])} \leq C |(I - \Delta)_{y^i}^{-1/2}\partial_{y^i} (\nabla_{x'} v_k |\nabla_{x'} \eta) |_{L_p(\mathbb{R}^n)}
\]
\[
\leq C |(\nabla_{x'} v_k |\nabla_{x'} \eta) |_{L_p(\mathbb{R}^n)}
\]
\[
\leq C|\nabla_{x'} \eta|_{\infty} |v_k|_{W^1_p(\mathbb{R}^n)}.
\]
In the same way we obtain
\[
|\partial_{y^i}^2 v_k |\nabla_{x'} \eta|^2 |_{W^{-1}_{p-1}([\mathbb{R}^n])} \leq C|\nabla_{x'} \eta|^2 |v_k|_{W^1_p(\mathbb{R}^n)},
\]
wheras
\[
|\partial_{y^i} v_k \Delta_{x'} \eta|_{W^{-1}_{p-1}([\mathbb{R}^n])} \leq C|v_k|_{L_p(\mathbb{R}^n)},
\]
since $\eta$ is smooth. Concerning the terms in the Neumann transmission condition, we obtain by trace theory
\[
|([\nabla_{x'} v_k]|\nabla_{x'} \eta)|_{W^{-1/p}_{p-1}([\mathbb{R}^{n-1}])} \leq C|\nabla_{x'} \eta|_{C^0([\mathbb{R}^{n-1}])} |\nabla_{x'} v_k|_{W^{-1/p}_{p-1}([\mathbb{R}^{n-1}])}
\]
\[
\leq C|\nabla_{x'} \eta|_{C^0([\mathbb{R}^{n-1}])} |v_k|_{W^{1-1/p}_{p-1}([\mathbb{R}^{n-1}])}
\]
\[
\leq C|\nabla_{x'} \eta|_{C^0([\mathbb{R}^{n-1}])} |v_k|_{W^2_p(\mathbb{R}^n)},
\]
where $\alpha \in (1/p, 1)$. These estimates show that the right hand side of (8.16) may be estimated by terms that are either of lower order or of highest order, but the higher order terms carry a factor of the form $|\nabla_{x'} \eta|_q^n$, $\theta > 0$, which becomes small, by decreasing the size of the chart $U_k$. Applying perturbation theory it follows that
there exists $\lambda_0 \geq 1$ such that for each chart $U_k$, $k = 2, \ldots, N$, the linear problem (8.10) has a bounded solution operator

$$S_k : W_p^{-1}(\mathbb{R}^n) \times W_p^{-1/p}([0,1]) \times W_p^{-1/p}([0,1]) \to W_p^{1}([0,1]),$$

provided $\lambda \geq \lambda_0$. This in turn yields that $\Theta^{-1}S_k\Theta$ is the corresponding solution operator for problem (8.15), i.e. we have

$$q_k = (\Theta^{-1}S_k\Theta)(F_k,G_k,H_k),$$

for each $k = 2, \ldots, N$, where $(F_k,G_k,H_k)$ denotes the right hand side of (8.15). Since $\Theta$ is a $C^2$-diffeomorphism, we obtain the estimate

$$\lambda^{1/2}\|q_k\|_{L_p(\Omega)} + |q|_{L_p(\Omega)} \leq C \left( |g|_{L_p^{-1}([0,1])} + |q|_{L_p(\Omega)} \right),$$

(8.17)

for some $s \in (1/p,1)$ and for each $k = 2, \ldots, N$. Here the constant $C > 0$ does not depend on $\lambda \geq \lambda_0$, as we have already shown in the investigation of (8.3). Let us introduce

$$|v|_{W_p^1(\Omega)} := |\lambda|^{1/2} |v|_{L_p(\Omega)} + |\nabla v|_{L_p(\Omega)}, \quad \lambda \geq 1, \quad v \in W_p^1(\Omega \setminus \Gamma),$$

which is an equivalent norm in $W_p^1(\Omega \setminus \Gamma)$. This yields

$$|g|_{W_p^1(\Omega)} \leq \sum_{k=0}^N |q_k|_{\lambda, W_p^1(\Omega)} \leq C \left( |g|_{L_p^{-1}([0,1])} + |q|_{L_p(\Omega)} \right),$$

with constants $C,M > 0$, being independent of $\lambda$. Since $s \in (1/p,1)$ we may apply interpolation theory to the result

$$|g|_{W_p^1(\Omega)} \leq \varepsilon |g|_{W_p^1(\Omega)} + C(\varepsilon)|g|_{L_p(\Omega)}$$

$$\leq \varepsilon |\lambda|_{W_p^1(\Omega)} + C(\varepsilon)|g|_{L_p(\Omega)}$$

$$\leq \left( \varepsilon + C(\varepsilon)/|\lambda|^{1/2} \right) |g|_{\lambda, W_p^1(\Omega)},$$

since by assumption $\lambda \geq 1$. Choosing first $\varepsilon > 0$ small enough and then $\lambda \geq 1$ sufficiently large, we finally obtain the estimate

$$|g|_{W_p^1(\Omega)} \leq C|g|_{W_p^{-1}([0,1])}$$

(8.18)

for the strong solution $g \in W_p^2(\Omega \setminus \Gamma)$ of (8.3). Now we want to reduce the regularity of $g$. Fix $g \in W_p^{-1}([0,1])$. Then there exists a sequence $(g_m) \subset W_p^{2-1/p}([0,1])$, such that $g_m \to g$ as $m \to \infty$ in $W_p^{2-1/p}([0,1])$. We denote by $q_m \in W_p^2(\Omega \setminus \Gamma)$ the corresponding solutions of (8.9). Then it follows from (8.18) that $(q_m)$ is a Cauchy sequence in $W_p^1(\Omega \setminus \Gamma)$. Therefore the limit $\lim_{m \to \infty} q_m =: q_\infty$ exists and $q_\infty \in W_p^1(\Omega \setminus \Gamma)$ is the unique weak solution of (8.9) for sufficiently large $\lambda \geq 1$.

**Lemma 8.3.** Let $1 < p < \infty$, $1/p + 1/p' = 1$ and let $g \in W_p^{-1}([0,1])$ be given. Then there exists $\lambda_0 \geq 1$ such that the problem

$$\lambda(q(\phi)_{L_2(\Omega)} + (\nabla q|\nabla \phi)_{L_2(\Omega)} = 0, \quad \phi \in H_p^1(\Omega), \quad [pq] = g, \quad x \in \Gamma,$$

has a unique solution $q \in H_p^1(\Omega \setminus \Gamma)$, provided $\lambda \geq \lambda_0$. Moreover, the solution $q \in H_p^1(\Omega \setminus \Gamma)$ satisfies the estimate

$$|q|_{H_p^1(\Omega \setminus \Gamma)} \leq C|g|_{W_p^{-1}([0,1])},$$

(8.19)
In a next step we consider the problem
\[
\lambda(q|\phi)_{L^2(\Omega)} + (\nabla q|\nabla \phi)_{L^2(\Omega)} = (f|\nabla \phi)_{L^2(\Omega)}, \quad \phi \in H^1_p(\Omega),
\]
\[
\rho q = 0, \quad x \in \Gamma,
\] (8.20)
where \( f \in L^p(\Omega; \mathbb{R}^n) \) is given. Observe that the mapping \( \psi_f : H^1_p(\Omega \setminus \Gamma) \to \mathbb{R} \) defined by
\[
\psi_f(\phi) := \langle \psi_f, \phi \rangle := \int_{\Omega}(f|\nabla \phi)dx,
\]
is linear and continuous, since
\[
|\psi_f(\phi)| \leq |f|_{L^p(\Omega; \mathbb{R}^n)}|\phi|_{H^1_p(\Omega)},
\]
hence \( \psi_f \in (H^1_p(\Omega \setminus \Gamma))^* \). With the help of the Dirichlet form
\[
a : H^1_p(\Omega \setminus \Gamma) \times H^1_p(\Omega \setminus \Gamma) \to \mathbb{R}, \quad a(q,v) := \int_{\Omega} \nabla q \cdot \nabla vdx,
\]
we define an operator \( A : H^1_p(\Omega \setminus \Gamma) \to (H^1_p(\Omega \setminus \Gamma))^* \) by means of
\[
\langle Aq, v \rangle := a(q,v),
\]
with domain \( D(A) = \{q \in H^1_p(\Omega \setminus \Gamma) : [\rho q] = 0 \text{ on } \Gamma\} \). Making use of these definitions, we may rewrite (8.20) in the abstract form
\[
\lambda q + Aq = \psi_f, \quad \text{in } (H^1_p(\Omega \setminus \Gamma))^*.
\] (8.21)
Since
\[
H^1_p(\Omega \setminus \Gamma) \hookrightarrow (H^1_p(\Omega \setminus \Gamma))^*,
\]
the resolvent of \( A \) is compact and therefore the spectral set \( \sigma(A) \) consists solely of a countably infinite sequence of isolated eigenvalues. By a bootstrap argument it is easily seen that the corresponding eigenfunctions are smooth. Hence, defining \( A_2 \) to be the part of \( A \) in \( L^2(\Omega, \Gamma) \) with domain \( D(A_2) = \{q \in D(A) : Aq \in L^2(\Omega, \Gamma)\} \), it follows that \( \sigma(A) = \sigma(A_2) \). Integrating by parts, we obtain
\[
D(A_2) = \{q \in H^2_p(\Omega \setminus \Gamma) : [\rho q] = 0, \quad [\partial_\nu q] = 0 \text{ on } \Gamma, \quad \partial_\nu q = 0 \text{ on } \partial \Omega\}
\]
and \( A_2q = -\Delta q \) in \( \Omega \setminus \Gamma \). Let \( \lambda \in \sigma(-A) = \sigma(-A_2) \) and let \( q \in D(A_2) \) be a corresponding eigenfunction. Then \( q \) satisfies the problem
\[
\lambda q - \Delta q = 0, \quad x \in \Omega \setminus \Gamma, \quad [\rho q] = 0, \quad x \in \Gamma, \quad [\partial_\nu q] = 0, \quad x \in \Gamma, \quad \partial_\nu q = 0, \quad x \in \partial \Omega.
\] (8.22)
Multiplying (8.22) by \( \rho q \) and integrating by parts, we obtain by (8.22) \( 2, 3, 4 \)
\[
-\lambda \int_{\Omega \setminus \Gamma} \rho(q^2)dx = -\int_{\Omega \setminus \Gamma} \rho \Delta qdx = -\rho_1 \int_{\Omega_1} q_1 \Delta q_1 dx - \rho_2 \int_{\Omega_2} q_2 \Delta q_2 dx
\]
\[
= \rho_1 \|
\frac{\partial q_1}{\partial x_1}\|^2 + \rho_2 \|
\frac{\partial q_2}{\partial x_2}\|^2 + \int_I (\partial_\nu q_2 \rho_2 q_2 - \partial_\nu q_1 \rho_1 q_1) d\Gamma
\]
\[
= \rho_1 \|
\frac{\partial q_1}{\partial x_1}\|^2 + \rho_2 \|
\frac{\partial q_2}{\partial x_2}\|^2 + \int_I \partial_\nu q_2 [\rho q] d\Gamma
\]
\[
= \rho_1 \|
\frac{\partial q_1}{\partial x_1}\|^2 + \rho_2 \|
\frac{\partial q_2}{\partial x_2}\|^2 \geq 0,
\]
where $q_j$ denotes the part of $q$ in $\Omega_j$. In particular it follows that $\lambda$ is real and $\lambda \leq 0$ for all $\lambda \in \sigma(-A)$ and if $\lambda = 0$ then $q_1$ and $q_2$ are both equal to a constant in $\Omega_1$ and $\Omega_2$, respectively, satisfying the identity $\rho_1 q_1 = \rho_2 q_2$. In other words the eigenvalue $\lambda = 0$ is simple and the kernel $N(A) = N(A_2)$ is given by

$$N(A) = K_{1, \rho}, \quad 1_{\rho}(x) := \chi_{\Omega_1}(x) + \frac{\rho_1}{\rho_2} \chi_{\Omega_2}(x), \quad x \in \Omega.$$  

Therefore, spectral theory implies $(H^1_p(\Omega \setminus \Gamma))^* = N(A) \oplus R(A)$ and $H^1_p(\Omega \setminus \Gamma) = N(A) \oplus Y$, where $Y$ is a closed subspace of $H^1_2(\Omega \setminus \Gamma)$. Note that these decompositions reduce the linear operator $A$. It follows that the equation $AQ = F$ has a unique solution $q \in Y \subset H^1_p(\Omega \setminus \Gamma)$ if and only if $F \in R(A)$, or equivalently $\langle F, 1_{\rho} \rangle = 0$. If $c \in \mathbb{K}$, then any other solution $\tilde{q} \in H^1_p(\Omega \setminus \Gamma)$ of $A\tilde{q} = F$ is given by $\tilde{q} = q + c1_{\rho}$ and we have the estimate

$$|\nabla \tilde{q}|_{L^p(\Omega)} \leq C|F|_{(H^1_p(\Omega \setminus \Gamma))^*}.$$

**Lemma 8.4.** Let $1 < p < \infty$, $1/p + 1/p' = 1$ and let $f \in L^p(\Omega; \mathbb{R}^n)$ be given. Then the problem

$$(\nabla q | \nabla \phi)_{L^2(\Omega)} = (f | \nabla \phi)_{L^2(\Omega)}, \quad \phi \in H^1_{p'}(\Omega),$$

has a unique solution $q \in \dot{H}^1_p(\Omega \setminus \Gamma)$, satisfying the estimate

$$|\nabla q|_{L^p(\Omega)} \leq C|f|_{L^p(\Omega; \mathbb{R}^n)}.$$

For the final step, let $v \in H^1_p(\Omega \setminus \Gamma)$ be the unique solution of

$$\lambda_0 (v | \phi)_{L^2(\Omega)} + (\nabla v | \nabla \phi)_{L^2(\Omega)} = 0, \quad \phi \in H^1_{p'}(\Omega),$$

$$\|\rho v\| = g, \quad x \in \Gamma,$$

which is well-defined, thanks to Lemma 8.3. With the help of this solution $v$, we define a functional $\psi_v \in (H^1_{p'}(\Omega \setminus \Gamma))^*$ by

$$\psi_v(\phi) := \int_{\Omega} \nabla v \cdot \nabla \phi \, dx.$$  

By definition it holds that $\psi_v(1_{\rho}) = 0$. Since also $\psi_f(1_{\rho}) = 0$ for all $f \in L^p(\Omega; \mathbb{R}^n)$, Lemma 8.4 yields a unique solution $w \in \dot{H}^1_p(\Omega \setminus \Gamma)$ of

$$(\nabla w | \nabla \phi)_{L^2(\Omega)} = \psi_f(\phi) - \psi_v(\phi), \quad \phi \in H^1_{p'}(\Omega),$$

$$\|\rho w\| = 0, \quad x \in \Gamma.$$  

Finally, the sum $q := v + w \in \dot{H}^1_p(\Omega \setminus \Gamma)$ is the unique solution of

$$(\nabla q | \nabla \phi)_{L^2(\Omega)} = \psi_f(\phi) = (f | \nabla \phi)_{L^2(\Omega)}, \quad \phi \in H^1_{p'}(\Omega),$$

$$\|\rho q\| = g, \quad x \in \Gamma$$  

and we have the estimate

$$|\nabla q|_{L^p(\Omega)} \leq C \left( |f|_{L^p(\Omega; \mathbb{R}^n)} + |g|_{W^{1-1/p}_{p'}(\Gamma)} \right).$$
Theorem 8.5. Let $1 < p < \infty$, $1/p + 1/p' = 1$, $f \in L_p(\Omega; \mathbb{R}^n)$ and $q \in W_p^{1-1/p}(\Gamma)$ be given. Then the problem

$$ \langle \nabla q | \nabla \phi \rangle_{L^2(\Omega)} = \langle f | \nabla \phi \rangle_{L^2(\Omega)}, \quad \phi \in H_p^1(\Omega), $$

$$ \| \rho q \| = g, \quad x \in \Gamma $$

has a unique solution $u \in H_p^1(\Omega; \Gamma)$ satisfying the estimate

$$ \| \nabla q \|_{L^p(\Omega)} \leq C_1 \left( \| f \|_{L^p(\delta \partial \Omega; \mathbb{R}^n)} + \| g \|_{W_p^{1-1/p}(\delta \partial \Omega)} \right). $$

If $J = [0, a]$, $f = f(t, x)$, $f \in H_p^1(J; L_p(\Omega; \mathbb{R}^n))$, $g = 0$, then $q \in H_p^1(J; H_p^1(\Omega \setminus \Gamma))$ and

$$ \| \nabla q \|_{H_p^1(J; L_p(\Omega; \mathbb{R}^n))} \leq C_2 \| f \|_{H_p^1(J; L_p(\Omega; \mathbb{R}^n))}. $$

8.2. Higher regularity in the bulk phases. The next problem we consider, is about higher regularity in the bulk phases $\Omega \setminus \Gamma$. To be precise, we study the elliptic transmission problem

$$ \lambda q - \Delta q = f, \quad x \in \Omega \setminus \Gamma, $$

$$ \| \rho q \| = 0, \quad x \in \Gamma, $$

$$ \| \partial_\nu q \| = 0, \quad x \in \Gamma, $$

$$ \delta \partial_\nu q + (1 - \delta) q = 0, \quad x \in \partial \Omega, \quad \delta \in \{0, 1\}, $$

where $f \in L_p(\Omega) \cap W_p^s(\Omega \setminus \Gamma)$, $s > 0$, is given and $\lambda \geq 1$. It is our aim to find a unique solution $q \in W_p^{2+s}(\Omega \setminus \Gamma)$ of (8.23). Note that by Theorem 8.1 there exists a unique solution $q \in H_p^2(\Omega \setminus \Gamma)$ of (8.23). Moreover, there exists a constant $C > 0$ being independent of $\lambda \geq 1$ such that the estimate

$$ \| q \|_{H_p^2(\Omega \setminus \Gamma)} \leq C \| f \|_{L_p(\Omega)} $$

is valid. Thus, it remains to show that in addition $q \in W_p^{2+s}(\Omega \setminus \Gamma)$, provided $f \in L_p(\Omega) \cap W_p^s(\Omega \setminus \Gamma)$. For this purpose let $\partial \Omega \in C^3$ and cover the compact set $\Omega$ by a union of finitely many open sets $U_k$, $k = 0, \ldots, N$ which are subject to the following conditions

- $\partial \Omega \subset U_0$ and $U_0 \cap \Gamma = \emptyset$;
- $U_1 \subset \Omega_1$ and $U_1 \cap \Gamma = \emptyset$;
- $U_k \cap \Gamma \neq \emptyset, U_k \cap \partial \Omega = \emptyset$ for $k = 2, \ldots, N$ and

$$ \bigcup_{k=2}^N U_k \supset \Gamma. $$

For $k \geq 2$, the sets $U_k$ may be balls with a fixed but arbitrarily small radius $r > 0$. As before, let $\{ \varphi_k \}_{k=0}^N$ be a partition of unity, such that $\text{supp} \varphi_k \subset U_k$ and $0 \leq \varphi_k(x) \leq 1$ for all $x \in \Omega$. Let $q_k := q \varphi_k$ and $f_k := f \varphi_k$.

Multiplying (8.23) by $\varphi_0$ yields the problem

$$ \lambda q_0 - \Delta q_0 = f_0 - 2(\nabla q | \nabla \varphi_0) - \varphi_0, \quad x \in \Omega, $$

$$ \delta \partial_\nu q_0 + (1 - \delta) q_0 = \delta q_0 \varphi_0, \quad x \in \partial \Omega, \quad \delta \in \{0, 1\}. $$

Since $\varphi_0$ is smooth and $q \in H_p^2(\Omega)$, the right hand side $(F_0, G_0)$ in (8.25) is in $W_p^s(\Omega) \times W_p^{s-1/p}(\partial \Omega)$, at least for $s \in (0, 1)$. It follows from [38] Theorem 5.5.1
with given functions $F$ and $F$ remove the inhomogeneity may use the transformation described above, to reduce (8.27) to the problem with some function $\psi$, which we denote by $F_1$. Then the solution of (8.26) is given by

$$q_1 = (\lambda - \Delta)^{-1} F_1.$$  

If $\alpha \in \{0, 1\}$ and $F_1 \in H^{\alpha}_{\mu}(\mathbb{R}^n)$ then $q_1 \in H^{2+\alpha}_{\mu}(\mathbb{R}^n)$ and

$$|q_1|_{H^{2+\alpha}_{\mu}(\mathbb{R}^n)} \leq C|I - \Delta|^{1+\alpha/2} F_1|_{L^{\infty}(\mathbb{R}^n)} \leq C|I - \Delta|^{1+\alpha/2} (\lambda - \Delta)^{-1} F_1|_{L^{\infty}(\mathbb{R}^n)}$$

$$= C|I - \Delta|C|I - \Delta|^{1+\alpha/2} F_1|_{L^{\infty}(\mathbb{R}^n)} \leq C||I - \Delta|C|I - \Delta|^{1+\alpha/2} F_1|_{L^{\infty}(\mathbb{R}^n)}$$

since $|I - \Delta|^{1+\alpha/2} \cdot |L^{\infty}(\mathbb{R}^n), \alpha \in \{0, 1\}$. Note that the term

$$||I - \Delta|C|I - \Delta|^{1+\alpha/2} F_1|_{L^{\infty}(\mathbb{R}^n)}$$

is independent of $\lambda \geq 1$, which follows e.g. from functional calculus. The real interpolation method and (8.27) then yield the estimate

$$|q_1|_{W^{2+\alpha}_{\mu}(\mathbb{R}^n)} \leq C|F_1|_{W^{\alpha}_{\mu}(\mathbb{R}^n)} \leq C \left( |F|_{W^{\alpha}_{\mu}(\mathbb{R}^n)} + |F|_{L^{\infty}(\mathbb{R}^n)} \right),$$

for $s \in (0, 1]$, where $C > 0$ does not depend on $\lambda \geq 1$. Next, we multiply (8.28) by $\varphi_k$, $k \in \{2, \ldots, N\}$, to obtain the pure transmission problems

$$\lambda \varphi_k - \Delta \varphi_k = f_k - 2(\nabla q|\nabla \varphi_k) - q\Delta \varphi_k, \quad x \in \mathbb{R}^n \setminus \Gamma,$$

$$[\rho \varphi_k] = 0, \quad x \in \Gamma,$$

$$[\partial_\nu \varphi_k] = [\varphi] \partial_\nu \varphi_k, \quad x \in \Gamma,$$

with some function $f_k \in W^{s}_{\mu}(\mathbb{R}^n \setminus \Gamma) \cap L^p(\mathbb{R}^n)$. For each fixed $k \in \{2, \ldots, N\}$ we may use the transformation described above, to reduce (8.27) to the problem

$$\lambda \psi - \Delta \psi = F, \quad x' \in \mathbb{R}^{n-1}, \quad y \in \mathbb{R},$$

$$[\rho \psi] = 0, \quad x' \in \mathbb{R}^{n-1}, \quad y = 0,$$

$$[\partial_\nu \psi] = G, \quad x' \in \mathbb{R}^{n-1}, \quad y = 0,$$

with given functions $F \in W^{s}_{\mu}(\mathbb{R}^n)$ and $G \in W^{1+s-1/p}_{\mu}([\mathbb{R}^{n-1} \setminus \Gamma]) \cap L^p(\mathbb{R}^n)$. First we remove the inhomogeneity $F$. To this end we solve the Dirichlet problems

$$\lambda \psi - \Delta \psi^+ - \partial_\nu^2 \psi^+ = F^+, \quad x' \in \mathbb{R}^{n-1}, \quad y > 0, \quad \psi^+(x', 0) = 0,$$

and

$$\lambda \psi - \Delta \psi^- - \partial_\nu^2 \psi^- = F^-, \quad x' \in \mathbb{R}^{n-1}, \quad y < 0, \quad \psi^-(x', 0) = 0,$$

where $F^+ := F|_{y>0}$ and $F^- := F|_{y<0}$. Let $\psi^\pm \in W^{2+s}_{\mu}(\mathbb{R}^n)$ be defined as

$$\psi^\pm(x', y) := \begin{cases} 
\psi^+(x', y), & y > 0, \\
\psi^-(x', y), & y < 0.
\end{cases}$$
Since $\psi^+(x', 0) = \psi^-(x', 0) = 0$ and $[\rho \tilde{\psi}] = [\rho] \psi^\pm = 0$, the shifted function $\tilde{\psi} := \psi - \psi^\pm$ solves the problem

$$
\lambda \tilde{\psi} - \Delta \tilde{\psi} = 0, \quad x' \in \mathbb{R}^{n-1}, \ y \in \mathbb{R},
$$

$$
[\partial_y \tilde{\psi}] = 0, \quad \tilde{\psi} = \tilde{\psi}^+, \quad x' \in \mathbb{R}^{n-1}, \ y = 0,
$$

$$
[\rho \tilde{\psi}] = 0, \quad \tilde{\psi} = \tilde{\psi}^-, \quad x' \in \mathbb{R}^{n-1}, \ y = 0,
$$

(8.29)

where $\tilde{G} := G - [\partial_y \psi^\pm] \in W_p^{1+s-1/p}(\mathbb{R}^{n-1})$. According to (8.6) and (8.7) the unique solution of (8.29) is given by

$$
\tilde{\psi}(y) = \frac{1}{\rho_1 + \rho_2} L^{-1} \left\{ \begin{array}{ll}
\rho_1 e^{-L y} \tilde{G}, & y > 0, \\
\rho_2 e^{L y} \tilde{G}, & y < 0,
\end{array} \right.
$$

(8.30)

where $L := (\lambda - \Delta_{x'})^{1/2}$ with domain $D(L) = H_p^1(\mathbb{R}^{n-1})$. Assume for a moment that $\tilde{G} \in W_p^{2-1/p}(\mathbb{R}^{n-1})$. Then it follows from semigroup theory and (8.30) that the solution of (8.29) satisfies the estimates

$$
|\tilde{\psi}|_{H_p^1(\mathbb{R}^n)} \leq C|\tilde{G}|_{W_p^{2-1/p}(\mathbb{R}^{n-1})}
$$

as well as

$$
|\tilde{\psi}|_{H_p^2(\mathbb{R}^n)} \leq C|\tilde{G}|_{W_p^{2-1/p}(\mathbb{R}^{n-1})},
$$

where the constant $C > 0$ does not depend on $\lambda \geq 1$. This can be seen as in the proof of Lemma 8.3. Applying the real interpolation method yields

$$
|\tilde{\psi}|_{W_p^{2+s}(\mathbb{R}^n)} \leq C|\tilde{G}|_{W_p^{1+s-1/p}(\mathbb{R}^{n-1})},
$$

for some $s \in (0, 1]$ and if $\tilde{G} \in W_p^{1+s-1/p}(\mathbb{R}^{n-1})$. We have thus shown that the transmission problem (8.28) has a unique solution $\psi \in W_p^{2+s}(\mathbb{R}^n)$ if and only if $F \in W_p^s(\mathbb{R}^n)$ and $G \in W_p^{1+s-1/p}(\mathbb{R}^{n-1})$. By perturbation theory, there exists $\lambda_0 \geq 1$ such that (8.27) has a unique solution $q_k \in W_p^{2+s}(\mathbb{R}^n \setminus \Gamma)$, $s \in (0, 1]$, satisfying the estimate

$$
|q_k|_{W_p^{2+s}(\mathbb{R}^n \setminus \Gamma)} \leq C \left( |f_k|_{W_p^s(\Omega, \Gamma)} + |(\nabla q) \nabla \varphi_k|_{W_p^s(\Omega, \Gamma)} + |\eta \Delta \varphi_k|_{W_p^s(\Omega, \Gamma)} \right)
$$

$$
+ |[q] \partial_{\nu} \varphi_k|_{W_p^{1+s-1/p}(\Gamma)},
$$

provided $\lambda \geq \lambda_0$. By the smoothness of $\varphi_k$ and by (8.24) we obtain the estimate

$$
|q_k|_{W_p^{2+s}(\mathbb{R}^n \setminus \Gamma)} \leq C \left( |f|_{W_p^s(\Omega, \Gamma)} + |g|_{W_p^{1+s}(\Omega, \Gamma)} \right) \leq C \left( |f|_{W_p^s(\Omega, \Gamma)} + |f|_{L_p(\Gamma)} \right),
$$

valid for all $k \in \{2, \ldots, N\}$ and $s \in (0, 1]$. Since $\{\varphi_k\}_{k=0}^N$ is a partition of unity, we obtain

$$
|q|_{W_p^{2+s}(\Omega, \Gamma)} \leq \sum_{k=0}^N |q_k|_{W_p^{2+s}(\Omega, \Gamma)} \leq C \left( |f|_{W_p^s(\Omega, \Gamma)} + |f|_{L_p(\Gamma)} \right),
$$

showing that $q \in W_p^{2+s}(\Omega, \Gamma)$, $s \in (0, 1]$. It is easy to extend this result to the case $\lambda \in [0, \lambda_0)$. To this end, let $f \in L_p(\Omega) \cap W_p^s(\Omega, \Gamma) \cap R(A_\delta)$, $s > 0$, where $A_\delta : H_p^2(\Omega, \Gamma) \to L_p(\Omega)$ was defined at the beginning of Section 3. Note that $R(A_\delta) = \{ f \in L_p(\Omega) : |f|^2_{L_p} = 0 \}$ if $\delta = 1$ and $\lambda = 0$ and $R(A_\delta) = L_p(\Omega)$ if either $\delta = 0$ and $\lambda \geq 0$ or $\delta = 1$ and $\lambda > 0$. Consider the solution $q \in H_p^2(\Omega, \Gamma)$ of (8.29).
with \( \lambda \in [0, \lambda_0) \), which is well-defined thanks to Theorem \( \text{S.1} \) and which satisfies the estimate (8.24). Rewriting (8.23) as

\[ \lambda_0 q - \Delta q = f + (\lambda_0 - \lambda)q, \]

we may regard the new right hand side \( f + (\lambda_0 - \lambda)q \) as a given function, say \( \tilde{f} \in W^{2+\delta}_p(\Omega \setminus \Gamma) \), \( s \in (0, 1] \). The above result for (8.23) then yields the estimate

\[
|q|_{W^{2+\delta}_p(\Omega \setminus \Gamma)} \leq C \left( |\tilde{f}|_{W^{2+\delta}_p(\Omega \setminus \Gamma)} + |\tilde{f}|_{L^p(\Omega)} \right)
\]

\[
\leq C \left( |f|_{W^{2+\delta}_p(\Omega \setminus \Gamma)} + |f|_{L^p(\Omega)} \right),
\]

since \( |q|_{W^{2+\delta}_p(\Omega \setminus \Gamma)} = |q|_{W^2_p(\Omega)} \leq C|q|_{H^2_p(\Omega)} \leq C|f|_{L^p(\Omega)} \),

by the smoothness of \( q \) and by (8.24). If \( s > 1 \) and \( f \in L^p(\Omega) \cap W^{s}_p(\Omega \setminus \Gamma) \), then \( q \in H^2_p(\Omega \setminus \Gamma) \), since \( f \in L^p(\Omega) \cap H^s_p(\Omega \setminus \Gamma) \). This additional regularity for \( q \) and the preceding steps allow us to conclude that \( q \in W^{2+\delta}_p(\Omega \setminus \Gamma) \), at least for \( s \in [1, 2) \).

By an obvious argument it follows that \( q \in W^{2+\delta}_p(\Omega \setminus \Gamma) \) for each fixed \( s > 0 \), provided \( f \in L^p(\Omega) \cap W^{s}_p(\Omega \setminus \Gamma) \). This yields the following result.

**Theorem 8.6.** Let \( \Omega \subset \mathbb{R}^n \) be a bounded domain with boundary \( \partial \Omega \subset C^{2+s} \), let \( 1 < p < \infty \), \( s > 0 \) and \( f \in L^p(\Omega) \cap W^{s}_p(\Omega \setminus \Gamma) \). Then the following assertions hold.

1. If \( \delta = 1 \) and \( \lambda = 0 \), then there exists a unique solution \( q \in W^{2+\delta}_p(\Omega \setminus \Gamma) \) of (8.23), provided that \( (f|1_p) = 0 \).

2. If either \( \delta = 1 \) and \( \lambda > 0 \) or \( \delta = 0 \) and \( \lambda \geq 0 \), then there exists a unique solution \( q \in W^{2+\delta}_p(\Omega \setminus \Gamma) \) of (8.23).

If in addition \( J = [0, a] \), \( f = f(t, x) \) and \( f \in H^1_p(J; L^p(\Omega) \cap W^{2+\delta}_p(\Omega \setminus \Gamma)) \) s.t. \( f(t, \cdot) \in R(A) \) for a.e. \( t \in J \), then \( q \in H^1_p(J; W^{2+\delta}_p(\Omega \setminus \Gamma) \cap N(A)) \).

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CENTER FOR SMART INTERFACES & INTERNATIONAL RESEARCH TRAINING GROUP "MATHEMATICAL FLUID DYNAMICS", TECHNICAL UNIVERSITY DARMSTADT, PETERSENSTR. 32, D-64287 DARMSTADT, GERMANY
E-mail address: koehne@csi.tu-darmstadt.de

INSTITUT FÜR MATHEMATIK, MARTIN-LUTHER-UNIVERSITÄT HALLE-WITTMENBERG, THEODOR-LIESER-STR. 5, D-06120 HALLE, GERMANY
E-mail address: jan.pruess@mathematik.uni-halle.de

INSTITUT FÜR MATHEMATIK, MARTIN-LUTHER-UNIVERSITÄT HALLE-WITTMENBERG, THEODOR-LIESER-STR. 5, D-06120 HALLE, GERMANY
E-mail address: mathias.wilke@mathematik.uni-halle.de