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# The Drag in a Navier-Stokes Flow with Friction-Driven Boundary Conditions 

Matthieu Bonnivard


#### Abstract

We consider the drag of an obstacle in a Navier-Stokes flow, associated to the friction-driven boundary conditions introduced by Bucur, Feireisl and Nečasová [6]. These boundary conditions account for the asymptotic effect of rough boundaries on the solutions to the Navier-Stokes equations. Consequently, they appear as a natural control variable in order to reduce the drag using the effect of micro-rugosities. In this paper, we prove the existence of a drag associated to friction-driven boundary conditions, we give a uniqueness criterion and, in the particular case of Navier's friction law, we prove that the drag is differentiable with respect to the friction coefficient, and compute its gradient.


## 1 Introduction

In this paper, we are interested in the definition and the mathematical properties of the drag of an obstacle in a Navier-Stokes flow, associated to the so-called friction-driven boundary conditions on the solid-fluid interface [6].

Let $N \in\{2,3\}$ be the space dimension. We define the obstacle $\mathcal{S}$ as a compact subset of $\mathbb{R}^{N}$, with a Lipschitz boundary. The friction-driven boundary conditions are determined by a triplet $\{\mu, A, \mathcal{V}\}$, where:

- $\mu$ is a capacitary measure concentrated on $\partial \mathcal{S}$,
- $A$ is a positive symmetric matrix function,
- $\mathcal{V}=\{\mathcal{V}(x)\}_{x \in \partial \mathcal{S}}$ is a collection of vector spaces, where for q.e. $x \in \partial \mathcal{S}, \mathcal{V}(x)$ is a subspace of the tangent hyperplane to $\mathcal{S}$ at $x$.

Let $\mathbf{u}_{\infty} \in \mathbb{R}^{N}$ be the velocity of the fluid at infinity. To define the drag of the solid $\mathcal{S}$ in the direction $\mathbf{u}_{\infty}$, we (artificially) introduce an open subset $\Omega \subset \mathbb{R}^{N}$ such that $\mathcal{S} \subset \Omega$, and define the fluid domain

$$
\mathcal{F}=\Omega \backslash \mathcal{S}
$$

We note $L_{0}^{2}(\mathcal{F})$ the subspace of $L^{2}(\mathcal{F})$ consisting of functions with mean value equal to zero. We consider the following problem: find $(\mathbf{u}, p) \in H^{1}\left(\mathcal{F}, \mathbb{R}^{N}\right) \times L_{0}^{2}(\mathcal{F})$ such that

$$
\begin{align*}
& -\operatorname{div}(\sigma(\mathbf{u}, p))+(\mathbf{u} \cdot \nabla) \mathbf{u}=0 \text { in } \mathcal{F}  \tag{1}\\
& \operatorname{div} \mathbf{u}=0 \text { in } \mathcal{F}  \tag{2}\\
& \mathbf{u}=\mathbf{u}_{\infty} \text { on } \partial \Omega  \tag{3}\\
& \mathbf{u}(x) \in \mathcal{V}(x) \text { for q.e. } x \in \partial \mathcal{S}  \tag{4}\\
& {[2 \nu \mathbf{D}(\mathbf{u}) \mathbf{n}+\mu A \mathbf{u}] \cdot \phi=0 \text { for any } \phi \in \mathcal{V}(x), x \in \partial \mathcal{S},} \tag{5}
\end{align*}
$$

where $\sigma(\mathbf{u}, p)$ is the stress tensor defined by

$$
\sigma(\mathbf{u}, p)=2 \nu \mathbf{D}(\mathbf{u})-p \mathrm{Id}
$$

$\nu>0$ being the viscosity of the fluid and $\mathbf{D}(\mathbf{u})$ the symmetric part of $\nabla \mathbf{u}$, defined by

$$
\mathbf{D}(\mathbf{u})=\frac{1}{2}\left((\nabla \mathbf{u})+(\nabla \mathbf{u})^{T}\right) .
$$

In condition (5), $\mathbf{n}$ stands for the normal vector field to $\partial \mathcal{S}$ oriented towards the solid.
Condition (4) expresses that the flow is driven by the family of linear spaces $\{\mathcal{V}(x)\}_{x \in \partial \mathcal{S}}$. Above, $[2 \nu \mathbf{D}(\mathbf{u}) \mathbf{n}+\mu A \mathbf{u}] \cdot \phi=0$ for $\phi \in \mathcal{V}(x)$ is a formal pointwise relation, which has to be understood globally on $\partial \mathcal{S}$. Precisely, condition (5) holds provided that

$$
2 \nu \int_{\partial \mathcal{S}} \mathbf{D}(\mathbf{u}) \mathbf{n} \cdot \phi \mathrm{d} x+\int_{\partial \mathcal{S}} A \mathbf{u} \cdot \phi \mathrm{~d} \mu=0
$$

for every $\phi \in H^{1}\left(\mathcal{F}, \mathbb{R}^{N}\right)$ such that $\phi(x) \in \mathcal{V}(x)$ for q.e. $x \in \partial S$. In particular, condition (5) is a friction condition, characterized by the measure $\mu$ and the matrix $A$.

Friction-driven boundary conditions have been obtained by Bucur, Feireisl and Nečasová in [6], as the general expression of the boundary conditions resulting from the rugosity effect on perfectly slippery walls. The rugosity effect has been investigated in [11, 10, 8, 2] and in the more recent papers $[7,6,5]$. It consists in approximating a given domain $U$ by a sequence of geometric perturbations $\left(U_{\varepsilon}\right)_{\varepsilon>0}$, and studying the asymptotic behaviour of a family of vector fields $\mathbf{u}_{\varepsilon} \in H^{1}\left(U_{\varepsilon}, \mathbb{R}^{N}\right)$ in the neighborhood of the boundary of the limit domain $U$. In [6], the authors have considered a sequence $\left(\mathbf{u}_{\varepsilon}\right)_{\varepsilon>0}$ of weak solutions to the NavierStokes equations with perfect slip boundary conditions in equi-Lipschitz domains $\left(U_{\varepsilon}\right)_{\varepsilon>0}$. They have proved the existence of a triplet $\{\mu, A, \mathcal{V}\}$, such that, up to a subsequence, $\left(\mathbf{u}_{\varepsilon}\right)_{\varepsilon>0}$ converges weakly to a function $\mathbf{u}$ solution to the Navier-Stokes equations with friction-driven boundary conditions.

This asymptotic process is an indirect proof of the existence of a solution to the NavierStokes equations with friction-driven boundary conditions, for admissible triplets $\{\mu, A, \mathcal{V}\}$, that is, triplets that can be effectively obtained by a sequence of geometric perturbations of the domain. However, since this result relies on abstract $\Gamma$-convergence tools, the class of admissible triplets cannot be explicitely described, at least for very general perturbations. The purpose of this paper is to prove that, in fact, problem (1)-(5) admits at least one weak solution (in a sense that will be precised below), for every $\nu>0, \mathbf{u}_{\infty} \in \mathbb{R}^{N}$ and every triplet
$\{\mu, A, \mathcal{V}\}$. This result is stated in Theorem 1. Moreover, we prove that for a suitable choice of $\nu, \mathbf{u}_{\infty}$ and $\mathcal{V}$, which does not depend on $\{\mu, A\}$, every weak solution to System (1)-(5) is unique (see Theorem 2).

In the second part of the paper, we study the differentiability of the drag function with respect to the friction part of the boundary condition. Let $\nu>0$ and $\mathbf{u}_{\infty} \in \mathbb{R}$. We can associate to every triplet $\{\mu, A, \mathcal{V}\}$ and every weak solution $\mathbf{u}$ to System (1)-(5), a drag $\mathcal{T}(\mu, A, \mathcal{V}, \mathbf{u})$ defined by

$$
\mathcal{T}(\mu, A, \mathcal{V}, \mathbf{u})=2 \nu \int_{\mathcal{F}}|\mathbf{D}(\mathbf{u})|^{2} \mathrm{~d} x+\int_{\partial \mathcal{S}} A \mathbf{u} \cdot \mathbf{u} \mathrm{~d} \mu
$$

In the case where $\mathbf{u}$ is unique, the drag depends only on the choice of the triplet $\{\mu, A, \mathcal{V}\}$, which appears as a natural control variable in drag minimization problems. As a first step in the study of the regularity of the drag with respect to $\{\mu, A, \mathcal{V}\}$, we consider Navier's boundary conditions on $\partial \mathcal{S}$, which are given by the non penetration condition

$$
\begin{equation*}
\mathbf{u} \cdot \mathbf{n}=0 \text { on } \partial \mathcal{S} \tag{6}
\end{equation*}
$$

and by the friction law

$$
\begin{equation*}
2 \nu[\mathbf{D}(\mathbf{u}) \mathbf{n}]_{t a n}+\beta \mathbf{u}=0 \text { on } \partial \mathcal{S} \tag{7}
\end{equation*}
$$

Here, $\mathbf{n}$ is the normal vector field to $\partial \mathcal{S}$ oriented towards the solid and $\beta \in L^{2}(\partial \mathcal{S})$ is a friction coefficient. For any vector field $\mathbf{v}$, defined on $\partial \mathcal{S}$, we note $[\mathbf{v}]_{t a n}$ its tangential part, defined by

$$
[\mathbf{v}]_{t a n}=\mathbf{v}-(\mathbf{v} \cdot \mathbf{n}) \mathbf{n}
$$

Boundary conditions (6)-(7) correspond to the case where $\mathcal{V}(x)$ is equal to the tangent hyperplane to $\partial \mathcal{S}$ on q.e. $x$ on the boundary, $A \equiv \operatorname{Id}$ and $\mu$ is absolutely continuous with respect to the surface measure on $\partial \mathcal{S}$ with a density $\beta \in L^{2}(\partial \mathcal{S})$. We prove that the drag is differentiable with respect to $\beta$ and compute the gradient of the drag in the $L^{2}(\partial \mathcal{S})$ topology (see Theorem 3 and Corollary 4).

The structure of the paper is the following. In Section 2, we introduce the functional framework associated with System (1)-(3) and boundary conditions (4)-(5) and (6)-(7), and state our main results. In Section 3, we recall some abstract material from [9] which we use in the proof of Theorems 1 and 2. Section 4 is devoted to the proof of Theorem 3 and Corollary 4.

## 2 Statement of the Main Results

We note

$$
\mathbf{H}^{1}(\mathcal{F})=H^{1}\left(\mathcal{F}, \mathbb{R}^{N}\right)
$$

To define the weak solutions to System (1)-(5), we introduce the following subsets of $\mathbf{H}^{1}(\mathcal{F})$ :

$$
K=\left\{\mathbf{v} \in \mathbf{H}^{1}(\mathcal{F}) \mid \operatorname{div} \mathbf{v}=0 \text { in } L^{2}(\mathcal{F}), \mathbf{v}(x) \in \mathcal{V}(x) \text { for q.e. } x \in \partial \mathcal{S}, \mathbf{v}=\mathbf{u}_{\infty} \text { on } \partial \Omega\right\}
$$

and

$$
V=\left\{\phi \in \mathbf{H}^{1}(\mathcal{F}) \mid \operatorname{div} \phi=0 \text { in } L^{2}(\mathcal{F}), \mathbf{v}(x) \in \mathcal{V}(x) \text { for q.e. } x \in \partial \mathcal{S}, \phi=0 \text { on } \partial \Omega\right\} .
$$

Let us fix the family of linear spaces $\mathcal{V}$ such that $V$ is non empty. In particular, $V$ is a separable Hilbert space for the $\mathbf{H}^{1}$ scalar product. We note $V^{\prime}$ its dual space for the $\mathbf{H}^{1}$ norm. Let $\mu$ be a capacitary measure concentrated on $\partial \mathcal{S}$ and $A$ a positive symmetric matrix function.

We say that $\mathbf{u} \in \mathbf{H}^{1}(\mathcal{F})$ is a weak solution to system (1)-(5) if $\mathbf{u} \in K$ and for every $\phi \in V$,

$$
\begin{equation*}
2 \nu \int_{\mathcal{F}} \mathbf{D}(\mathbf{u}): \mathbf{D}(\phi) \mathrm{d} x+\int_{\mathcal{F}}[(\mathbf{u} \cdot \nabla) \mathbf{u}] \cdot \phi \mathrm{d} x+\int_{\partial \mathcal{S}} A \mathbf{u} \cdot \phi \mathrm{~d} \mu=0 . \tag{8}
\end{equation*}
$$

The following results hold:
Theorem 1 Given $\nu>0$ and $\mathbf{u}_{\infty} \in \mathbb{R}^{N}$, there exists at least one weak solution to System (1)-(5). Moreover, there exists a constant $M>0$ depending only on $\nu, \mathbf{u}_{\infty}, \Omega$ and $\mathcal{S}$, such that every weak solution $\mathbf{u} \in K$ satisfies the following estimate:

$$
\begin{equation*}
\|\mathbf{u}\|_{\mathbf{H}^{1}(\mathcal{F})} \leqslant M . \tag{9}
\end{equation*}
$$

Theorem 2 There exists a constant $C=C(\Omega, \mathcal{S})>0$ such that, for all $\mathbf{u}_{\infty} \in \mathbb{R}^{N}$, there exists $\nu_{0}>0$ satisfying

$$
\begin{equation*}
\nu_{0} \leqslant C\left(\left|\mathbf{u}_{\infty}\right|+\left|\mathbf{u}_{\infty}\right|^{2}\right)^{1 / 2} \tag{10}
\end{equation*}
$$

and such that for all $\nu>\nu_{0}$, System (1)-(5) has a unique weak solution, which satisfies the following estimate:

$$
\begin{equation*}
\|\mathbf{u}\|_{\mathbf{H}^{1}(\mathcal{F})} \leqslant C^{\prime}\left(\nu+\left|\mathbf{u}_{\infty}\right|\right) \tag{11}
\end{equation*}
$$

where $C^{\prime}$ is a positive constant depending only on $\Omega, \mathcal{S}$.
Now we consider System (1)-(3) with Navier's boundary conditions (6)-(7) on $\partial \mathcal{S}$. To define weak solutions to this system, we introduce

$$
\mathbf{H}_{\tau}=\left\{\mathbf{v} \in \mathbf{H}^{1}(\mathcal{F}) \mid \operatorname{div} \mathbf{v}=0 \text { in } L^{2}(\mathcal{F}), \mathbf{v} \cdot \mathbf{n}=0 \text { on } \partial \mathcal{S}\right\}
$$

and define, as above,

$$
K_{\tau}=\left\{\mathbf{v} \in \mathbf{H}_{\tau} \mid \mathbf{v}=\mathbf{u}_{\infty} \text { on } \partial \Omega\right\}
$$

and

$$
V_{\tau}=\left\{\phi \in \mathbf{H}_{\tau} \mid \phi=0 \text { on } \partial \Omega\right\}
$$

We say that $\mathbf{u} \in \mathbf{H}^{1}(\mathcal{F})$ is a weak solution to system (1)-(3) with boundary conditions (6)-(7) if $\mathbf{u} \in K_{\tau}$ and for every $\phi \in V_{\tau}$,

$$
\begin{equation*}
2 \nu \int_{\mathcal{F}} \mathbf{D}(\mathbf{u}): \mathbf{D}(\phi) \mathrm{d} x+\int_{\mathcal{F}}[(\mathbf{u} \cdot \nabla) \mathbf{u}] \cdot \phi \mathrm{d} x+\int_{\partial \mathcal{S}} \beta \mathbf{u} \cdot \phi \mathrm{d} \mathcal{H}^{N-1}=0 . \tag{12}
\end{equation*}
$$

Let $\mathbf{u}_{\infty} \in \mathbb{R}^{N}$ and $\nu_{0}>0$ as in Theorem 2. We fix $\nu>\nu_{0}$, so that for every $\beta \in L^{2}(\partial \mathcal{S})$, any weak solution to system (1)-(3), (6)-(7), is unique. To each friction coefficient $\beta \in$ $L^{2}(\partial \mathcal{S})$, we associate a drag $T(\beta)$ defined by

$$
T(\beta)=2 \nu \int_{\mathcal{F}}|\mathbf{D}(\mathbf{u})|^{2} \mathrm{~d} x+\int_{\partial \mathcal{S}} \beta|\mathbf{u}|^{2} \mathrm{~d} \mathcal{H}^{N-1}
$$

where $\mathbf{u} \in K_{\tau}$ satisfies (12).
To address the differentiability of the drag with respect to $\beta$, we introduce the following open subset of $L^{2}(\partial \mathcal{S})$

$$
\mathcal{O}=\left\{\beta \in L^{2}(\partial \mathcal{S}) \mid \beta>0 \text { q.e. on } \partial \mathcal{S}\right\},
$$

and define the mapping

$$
\begin{aligned}
R: \mathcal{O} & \rightarrow \mathbf{H}_{\tau} \\
\beta & \mapsto \mathbf{u},
\end{aligned}
$$

where $\mathbf{u} \in K_{\tau}$ satisfies (12). The following result holds:
Theorem 3 The mapping

$$
\beta \in \mathcal{O} \rightarrow R(\beta) \in \mathbf{H}_{\tau}
$$

is differentiable, and for every $\beta \in \mathcal{O}$, its derivative $D R(\beta) \in \mathcal{L}\left(L^{2}(\partial \mathcal{S}), \mathbf{H}_{\tau}\right)$ at point $\beta$ is defined by

$$
D R(\beta) h=\mathbf{w} \quad \forall h \in L^{2}(\partial \mathcal{S})
$$

where $\mathbf{w}$ is the unique solution to the system:

$$
\left\{\begin{align*}
-\operatorname{div}(\sigma(\mathbf{w}, p))+(\mathbf{w} \cdot \nabla) \mathbf{u}+(\mathbf{u} \cdot \nabla) \mathbf{w}=0 & \text { in } \mathcal{F}  \tag{13}\\
\operatorname{div} \mathbf{w}=0 & \text { in } \mathcal{F} \\
\mathbf{w}=0 & \text { on } \partial \Omega \\
\mathbf{w} \cdot \mathbf{n}=0 & \text { on } \partial \mathcal{S} \\
2 \nu[\mathbf{D}(\mathbf{w}) \mathbf{n}]_{\tan }+\beta \mathbf{w}=-h \mathbf{u} & \text { on } \partial \mathcal{S}
\end{align*}\right.
$$

in the following sense: $\mathbf{w} \in V_{\tau}$ and for every $\phi \in V_{\tau}$,
$2 \nu \int_{\mathcal{F}} \mathbf{D}(\mathbf{w}): \mathbf{D}(\phi) \mathrm{d} x+\int_{\mathcal{F}}[(\mathbf{w} \cdot \nabla) \mathbf{u}+(\mathbf{u} \cdot \nabla) \mathbf{w}] \cdot \phi \mathrm{d} x+\int_{\partial \mathcal{S}} \beta \mathbf{w} \cdot \phi \mathrm{d} \mathcal{H}^{N-1}=-\int_{\partial \mathcal{S}} h \mathbf{u} \cdot \phi \mathrm{~d} \mathcal{H}^{N-1}$.

The following corollary has interesting applications to the drag optimization problem.
Corollary 4 The mapping

$$
\beta \in \mathcal{O} \rightarrow T(\beta) \in \mathbb{R}
$$

is differentiable, and for every $\beta \in \mathcal{O}$, its derivative $D T(\beta) \in \mathcal{L}\left(L^{2}(\partial \mathcal{S}), \mathbb{R}\right)$ is defined by

$$
D T(\beta) h=2\left(2 \nu \int_{\mathcal{F}} D(\mathbf{u}): D(\mathbf{w}) \mathrm{d} x+\int_{\partial \mathcal{S}} \beta \mathbf{u} \cdot \mathbf{w}\right)+\int_{\partial \mathcal{S}} h|\mathbf{u}|^{2} \mathrm{~d} \mathcal{H}^{N-1} \quad \forall h \in L^{2}(\partial \mathcal{S})
$$

where $\mathbf{u}=R(\beta)$ and $\mathbf{w}=D R(\beta) h$ is the solution to system (13). For every $\beta$, the gradient of $T$ at $\beta$ with respect to the $L^{2}(\partial \mathcal{S})$ topology is given by the following formula:

$$
\begin{equation*}
\nabla T(\beta)=(\mathbf{u}+\psi) \cdot \mathbf{u} \quad \text { q.e. on } \partial \mathcal{S}, \tag{15}
\end{equation*}
$$

where $\mathbf{u}=R(\beta)$ and $\psi$ is the unique solution to

$$
\left\{\begin{align*}
-\operatorname{div}(\sigma(\psi, p))+(\nabla \mathbf{u})^{T} \psi-(\mathbf{u} \cdot \nabla) \psi & =2(\mathbf{u} \cdot \nabla) \mathbf{u} \text { in } \mathcal{F},  \tag{16}\\
\operatorname{div} \psi & =0 \text { in } \mathcal{F}, \\
\psi & =0 \text { on } \partial \Omega \\
\psi \cdot \mathbf{n} & =0 \text { on } \partial \mathcal{S} \\
2 \nu[\mathbf{D}(\psi) \mathbf{n}]_{t a n}+\beta \psi & =0 \text { on } \partial \mathcal{S}
\end{align*}\right.
$$

in the sense that $\psi \in V_{\tau}$ and for every $\phi \in V_{\tau}$,

$$
\begin{equation*}
2 \nu \int_{\mathcal{F}} \mathbf{D}(\psi): \mathbf{D}(\phi) \mathrm{d} x+\int_{\mathcal{F}}\left[(\nabla \mathbf{u})^{T} \psi-(\mathbf{u} \cdot \nabla) \psi\right] \cdot \phi \mathrm{d} x+\int_{\partial \mathcal{S}} \beta \psi \cdot \phi \mathrm{d} \mathcal{H}^{N-1}=2 \int_{\mathcal{F}}[(\mathbf{u} \cdot \nabla) \mathbf{u}] \cdot \phi \mathrm{d} x \tag{17}
\end{equation*}
$$

## 3 Proof of Theorems 1 and 2

In order to prove Theorems 1 and 2, we follow the standard steps from the theory of stationary Navier-Stokes equations, which we adapt to the framework of friction-driven boundary conditions. These arguments are developed in [9], for instance.

We define the bilinear form $a_{0}: \mathbf{H}^{1}(\mathcal{F}) \times \mathbf{H}^{1}(\mathcal{F}) \rightarrow \mathbb{R}$ by

$$
\begin{equation*}
a_{0}\left(\mathbf{v}_{1}, \mathbf{v}_{2}\right)=2 \nu \int_{\mathcal{F}} \mathbf{D}\left(\mathbf{v}_{1}\right): \mathbf{D}\left(\mathbf{v}_{2}\right) \mathrm{d} x+\int_{\partial \mathcal{S}} A \mathbf{v}_{1} \cdot \mathbf{v}_{2} \mathrm{~d} \mu \quad \forall\left(\mathbf{v}_{1}, \mathbf{v}_{2}\right) \in \mathbf{H}^{1}(\mathcal{F}) \times \mathbf{H}^{1}(\mathcal{F}), \tag{18}
\end{equation*}
$$

and the trilinear form $a_{1}: \mathbf{H}^{1}(\mathcal{F}) \times \mathbf{H}^{1}(\mathcal{F}) \times \mathbf{H}^{1}(\mathcal{F}) \rightarrow \mathbb{R}$ by

$$
\begin{equation*}
a_{1}\left(\mathbf{v}, \mathbf{v}_{1}, \mathbf{v}_{2}\right)=\int_{\mathcal{F}}\left[(\mathbf{v} \cdot \nabla) \mathbf{v}_{1}\right] \cdot \mathbf{v}_{2} \mathrm{~d} x \quad \forall\left(\mathbf{v}, \mathbf{v}_{1}, \mathbf{v}_{2}\right) \in \mathbf{H}^{1}(\mathcal{F}) \times \mathbf{H}^{1}(\mathcal{F}) \times \mathbf{H}^{1}(\mathcal{F}) \tag{19}
\end{equation*}
$$

Let $\mathbf{u}_{0} \in K$ such that

$$
\begin{equation*}
\left\|\mathbf{u}_{0}\right\|_{\mathbf{H}^{1}(\mathcal{F})} \leqslant C_{0}\left|\mathbf{u}_{\infty}\right| \tag{20}
\end{equation*}
$$

where $C_{0}>0$ depends only on $\mathcal{F}$. The existence of $C_{0}$ and $\mathbf{u}_{0} \in K$ satisfying (20) was proved by Bogovskiŭ [3]. Setting $\widetilde{\mathbf{u}}=\mathbf{u}-\mathbf{u}_{0}$, we get that a function $\mathbf{u} \in K$ satisfies (12) if and only if $\widetilde{\mathbf{u}} \in V$ and satisfies the following identity:

$$
\begin{array}{r}
\forall \phi \in V \quad a_{0}(\widetilde{\mathbf{u}}, \phi)+a_{1}\left(\mathbf{u}_{0}, \widetilde{\mathbf{u}}, \phi\right)+a_{1}\left(\widetilde{\mathbf{u}}, \mathbf{u}_{0}, \phi\right)+a_{1}(\widetilde{\mathbf{u}}, \widetilde{\mathbf{u}}, \phi) \\
=-a_{0}\left(\mathbf{u}_{0}, \phi\right)-a_{1}\left(\mathbf{u}_{0}, \mathbf{u}_{0}, \phi\right) . \tag{21}
\end{array}
$$

We define the mapping $a: \mathbf{H}^{1}(\mathcal{F}) \times \mathbf{H}^{1}(\mathcal{F}) \times \mathbf{H}^{1}(\mathcal{F}) \rightarrow \mathbb{R}$ by

$$
\begin{align*}
& a\left(\mathbf{v}, \mathbf{v}_{1}, \mathbf{v}_{2}\right)=a_{0}\left(\mathbf{v}_{1}, \mathbf{v}_{2}\right)+a_{1}\left(\mathbf{u}_{0}, \mathbf{v}_{1}, \mathbf{v}_{2}\right)+a_{1}\left(\mathbf{v}_{1}, \mathbf{u}_{0}, \mathbf{v}_{2}\right)+a_{1}\left(\mathbf{v}, \mathbf{v}_{1}, \mathbf{v}_{2}\right) \\
& \forall\left(\mathbf{v}, \mathbf{v}_{1}, \mathbf{v}_{2}\right) \in \mathbf{H}^{1}(\mathcal{F}) \times \mathbf{H}^{1}(\mathcal{F}) \times \mathbf{H}^{1}(\mathcal{F}), \tag{22}
\end{align*}
$$

and the linear form $\mathbf{f}=\mathbf{f}\left(\mathbf{u}_{0}\right) \in V^{\prime}$ by

$$
\langle\mathbf{f}, \phi\rangle_{V^{\prime} \times V}=-a_{0}\left(\mathbf{u}_{0}, \phi\right)-a_{1}\left(\mathbf{u}_{0}, \mathbf{u}_{0}, \phi\right) \quad \forall \phi \in V .
$$

The problem we consider is equivalent to the following : find $\mathbf{u}_{0} \in K$ and $\widetilde{\mathbf{u}} \in V$ such that

$$
\begin{equation*}
a(\widetilde{\mathbf{u}}, \widetilde{\mathbf{u}}, \phi)=\left\langle\mathbf{f}\left(\mathbf{u}_{0}\right), \phi\right\rangle_{V^{\prime} \times V} \quad \forall \phi \in V \tag{23}
\end{equation*}
$$

We recall the following results from [9].
Theorem 5 Assume that the following hypotheses hold:
(i) for all $\mathbf{v} \in \mathbf{H}^{1}(\mathcal{F})$, the mapping $\left(\mathbf{v}_{1}, \mathbf{v}_{2}\right) \rightarrow a\left(\mathbf{v}, \mathbf{v}_{1}, \mathbf{v}_{2}\right)$ is continuous on $\mathbf{H}^{1}(\mathcal{F}) \times$ $\mathbf{H}^{1}(\mathcal{F})$;
(ii) there exists a constant $\alpha>0$ such that

$$
\begin{equation*}
a(\phi, \phi, \phi) \geqslant \alpha\|\phi\|_{\mathbf{H}^{1}(\mathcal{F})}^{2} \quad \forall \phi \in V ; \tag{24}
\end{equation*}
$$

(iii) the space $V$ is separable and, for all $\phi \in V$, the mapping

$$
\mathbf{v} \rightarrow a(\mathbf{v}, \mathbf{v}, \phi)
$$

is sequentially weakly continuous on $V$, i.e.,

$$
\begin{equation*}
\text { weak } \lim _{m \rightarrow \infty} \mathbf{v}_{m}=\mathbf{v} \text { in V implies } \lim _{m \rightarrow \infty} a\left(\mathbf{v}_{m}, \mathbf{v}_{m}, \phi\right)=a(\mathbf{v}, \mathbf{v}, \phi) \quad \forall \phi \in V . \tag{25}
\end{equation*}
$$

Then, problem (23) admits at least one solution.
Theorem 6 Moreover, assume that
(iv) the bilinear form $a(\mathbf{v}, \cdot, \cdot)$ is uniformly $V$-elliptic with respect to $\mathbf{v}$, i.e., there exists a constant $\alpha>0$ such that

$$
\begin{equation*}
a(\mathbf{v}, \phi, \phi) \geqslant \alpha\|\phi\|_{\mathbf{H}^{1}(\mathcal{F})}^{2} \quad \forall(\mathbf{v}, \phi) \in V \times V \tag{26}
\end{equation*}
$$

(v) there exists a continuous and monotonically increasing function $L: \mathbb{R}_{+} \rightarrow \mathbb{R}_{+}$such that for all $\eta>0$

$$
\begin{gather*}
\left|a\left(\mathbf{v}_{1}, \mathbf{v}_{1}, \mathbf{v}_{2}\right)-a\left(\mathbf{v}_{2}, \mathbf{v}_{1}, \mathbf{v}_{2}\right)\right| \leqslant L(\eta)\left\|\mathbf{v}_{1}\right\|_{\mathbf{H}^{1}(\mathcal{F})}\left\|\mathbf{v}_{2}\right\|_{\mathbf{H}^{1}(\mathcal{F})}\left\|\mathbf{v}_{1}-\mathbf{v}_{2}\right\|_{\mathbf{H}^{1}(\mathcal{F})} \\
\forall \mathbf{v}_{1}, \mathbf{v}_{2} \in V, \quad \forall \mathbf{v}_{1}, \mathbf{v}_{2} \in S_{\eta}=\left\{\mathbf{v} \in V \mid\|\mathbf{v}\|_{\mathbf{H}^{1}(\mathcal{F})} \leqslant \eta\right\} . \tag{27}
\end{gather*}
$$

Then, under the condition

$$
\begin{equation*}
\left[\|\mathbf{f}\|_{V^{\prime}} /\left(\alpha^{2}\right)\right] L\left(\|\mathbf{f}\|_{V^{\prime}} / \alpha\right)<1 \tag{28}
\end{equation*}
$$

the solution $\widetilde{\mathbf{u}} \in V$ to problem (23) is unique, and satisfies the following estimate:

$$
\begin{equation*}
\|\widetilde{\mathbf{u}}\|_{\mathbf{H}^{1}(\mathcal{F})} \leqslant\|\mathbf{f}\|_{V^{\prime}} / \alpha . \tag{29}
\end{equation*}
$$

Proof of Theorem 1. Let us fix $\mathbf{u}_{\infty} \in \mathbb{R}^{N}$ and $\nu>0$. We need to prove that the mapping $a$ defined by (22) satisfies the assumptions of Theorem 5 .

By Sobolev imbedding theorem (see the standard monograph by Adams [1]), $\mathbf{H}^{1}(\mathcal{F})$ is continuously imbedded in $L^{4}\left(\mathcal{F}, \mathbb{R}^{N}\right)$ for $N=2,3$. Therefore, by Hölder inequality, the mapping $a_{1}$ defined by (19) is continuous. Since $a_{0}$ defined by (18) is continuous, we obtain the continuity of $a$, and consequently, condition $(i)$ from Theorem 5 holds.

To prove ( $i i$ ), we prove the stronger condition (iv). We use the following Korn inequality, which is proved for instance in [4] (see the proof of Lemma VII.3.5, p. 329):

Lemma 7 There exists $C_{1}>0$ such that for all $\phi \in V$,

$$
\|\nabla \phi\|_{L^{2}\left(\mathcal{F}, \mathbb{R}^{N \times N}\right)} \leqslant C_{1}\|D(\phi)\|_{L^{2}\left(\mathcal{F}, \mathbb{R}^{N \times N}\right)} .
$$

Consequently, by Poincaré inequality, there exists $C_{2}>0$ such that for all $\phi \in V$,

$$
\begin{equation*}
\frac{1}{C_{2}}\|\phi\|_{\mathbf{H}^{1}(\mathcal{F})}^{2} \leqslant 2\|D(\phi)\|_{L^{2}\left(\mathcal{F}, \mathbb{R}^{N \times N}\right)}^{2} . \tag{30}
\end{equation*}
$$

Let $\mathbf{v}, \phi \in V$. We observe that

$$
a(\mathbf{v}, \phi, \phi)=a_{0}(\phi, \phi)+a_{1}\left(\phi, \mathbf{u}_{0}, \phi\right) .
$$

Indeed, $a_{1}\left(\mathbf{u}_{0}, \phi, \phi\right)=0$ and $a_{1}(\mathbf{v}, \phi, \phi)=0$, which comes from the following computation:

$$
\begin{align*}
a_{1}\left(\mathbf{u}_{0}, \phi, \phi\right) & =\frac{1}{2} \int_{\mathcal{F}}\left(\mathbf{u}_{0} \cdot \nabla\right)\left(|\phi|^{2}\right) \mathrm{d} x \\
& =\frac{1}{2}\left(\int_{\partial \mathcal{S}}|\phi|^{2} \mathbf{u}_{0} \cdot \mathbf{n} \mathrm{~d} \mathcal{H}^{N-1}+\int_{\partial \Omega}|\phi|^{2} \mathbf{u}_{0} \cdot \mathbf{n} \mathrm{~d} \mathcal{H}^{N-1}\right) . \tag{31}
\end{align*}
$$

Since $\phi \in V, \phi=0$ on $\partial \Omega$, and since $\mathbf{u}_{0} \in K, \mathbf{u}_{0} \cdot \mathbf{n}=0$ on $\partial \mathcal{S}$. Therefore, last boundary integrals vanish. Replacing $\mathbf{u}_{0}$ by $\mathbf{v}$ we obtain the desired conclusion.

For every $\mathbf{u}_{0} \in K$ satisfying (20), we note $\rho\left(\mathbf{u}_{0}\right)$ the best positive constant such that

$$
\begin{equation*}
\left|a_{1}\left(\phi, \mathbf{u}_{0}, \phi\right)\right| \leqslant \rho\left(\mathbf{u}_{0}\right)\|\phi\|_{\mathbf{H}^{1}(\mathcal{F})}^{2} \quad \forall \phi \in V \tag{32}
\end{equation*}
$$

We prove the following result:
Lemma 8 For every $\delta>0$, there exists $\mathbf{u}_{0} \in K$ satisfying (20) and such that $\rho\left(\mathbf{u}_{0}\right) \leqslant \delta$.
Proof. We consider a bounded extension operator

$$
P: \mathbf{H}^{1}(\mathcal{F}) \rightarrow H_{0}^{1}\left(\Omega, \mathbb{R}^{N}\right)
$$

Let $\delta>0$. By the proof of Lemma 2.3 in [9], chapter IV, there exists $\widetilde{\mathbf{u}}_{0} \in H^{1}\left(\Omega, \mathbb{R}^{N}\right)$ such that

$$
\left\{\begin{array}{rlr}
\operatorname{div} \widetilde{\mathbf{u}}_{0} & =0 & \text { in } \Omega \\
\widetilde{\mathbf{u}}_{0} & =\mathbf{u}_{\infty} & \text { on } \partial \Omega \\
\widetilde{\mathbf{u}}_{0} & =0 \quad \text { in a neighborhood of } \mathcal{S}
\end{array}\right.
$$

and such that for every $\varphi \in H_{0}^{1}\left(\Omega, \mathbb{R}^{N}\right)$,

$$
\begin{equation*}
\left|\int_{\Omega}\left[(\varphi \cdot \nabla) \widetilde{\mathbf{u}}_{0}\right] \cdot \varphi \mathrm{d} x\right| \leqslant \frac{\delta}{\|P\|}\|\varphi\|_{H^{1}\left(\Omega, \mathbb{R}^{N}\right)} . \tag{33}
\end{equation*}
$$

Let $\phi \in V$. By inequality (33), we obtain

$$
\begin{aligned}
\int_{\Omega}\left[((P \phi) \cdot \nabla) \widetilde{\mathbf{u}}_{0}\right] \cdot(P \phi) \mathrm{d} x & \leqslant \frac{\delta}{\|P\|}\|P \phi\|_{H^{1}\left(\Omega, \mathbb{R}^{N}\right)} \\
& \leqslant \delta\|\phi\|_{\mathbf{H}^{1}(\mathcal{F})} .
\end{aligned}
$$

Since $\widetilde{\mathbf{u}}_{0}$ vanishes on $\mathcal{S}$, observing that

$$
\int_{\Omega}\left[((P \phi) \cdot \nabla) \widetilde{\mathbf{u}}_{0}\right] \cdot(P \phi) \mathrm{d} x=\int_{\mathcal{F}}\left[(\phi \cdot \nabla) \widetilde{\mathbf{u}}_{0}\right] \cdot \phi \mathrm{d} x
$$

we obtain the desired result with $\mathbf{u}_{0}=\left(\widetilde{\mathbf{u}}_{0}\right)_{\mid \mathcal{F}}$.
Let $\mathbf{u}_{0} \in K$ such that $\rho\left(\mathbf{u}_{0}\right)<\nu / C_{2}$, where $C_{2}$ is defined by (30). By conditions (30) and (32), for every $\mathbf{v}, \phi \in V$,

$$
a(\mathbf{v}, \phi, \phi) \geqslant\left(\frac{\nu}{C_{2}}-\rho\left(\mathbf{u}_{0}\right)\right)\|\phi\|_{\mathbf{H}^{1}(\mathcal{F})}^{2}
$$

Consequently, condition (iv) holds with $\alpha=\nu / C_{2}-\rho\left(\mathbf{u}_{0}\right)$.
To prove condition (25), we only need to consider the nonlinear term $a_{1}\left(\mathbf{v}_{m}, \mathbf{v}_{m}, \phi\right)$. Integrating by part and using the incompressibility condition, we obtain

$$
a_{1}\left(\mathbf{v}_{m}, \mathbf{v}_{m}, \phi\right)=-\int_{\mathcal{F}}\left[\left(\mathbf{v}_{m} \cdot \nabla\right) \phi\right] \cdot \mathbf{v} \mathrm{d} x+\int_{\partial \mathcal{S} \cup \partial \Omega}\left(\mathbf{v}_{m} \cdot \phi\right)\left(\mathbf{v}_{m} \cdot \mathbf{n}\right) \mathrm{d} \mathcal{H}^{N-1}
$$

Since $\phi, \mathbf{v}_{m} \in V$, we conclude as above that

$$
a_{1}\left(\mathbf{v}_{m}, \mathbf{v}_{m}, \phi\right)=-\int_{\mathcal{F}}\left[\left(\mathbf{v}_{m} \cdot \nabla\right) \phi\right] \cdot \mathbf{v}_{m} \mathrm{~d} x
$$

Since $\mathbf{H}^{1}(\mathcal{F})$ is continuously imbedded in $L^{4}\left(\mathcal{F}, \mathbb{R}^{N}\right)$, we obtain

$$
\lim _{m \rightarrow \infty} a_{1}\left(\mathbf{v}_{m}, \mathbf{v}_{m}, \phi\right)=a_{1}(\mathbf{v}, \mathbf{v}, \phi)
$$

This proves condition (iii).
Now let $\mathbf{u} \in K$ be a weak solution to System (1)-(5). We define $\widetilde{\mathbf{u}}=\mathbf{u}-\mathbf{u}_{0}$. Applying relation (23) to $\phi=\widetilde{\mathbf{u}}$ and using inequality (26), we obtain

$$
\|\widetilde{\mathbf{u}}\|_{\mathbf{H}^{1}(\mathcal{F})} \leqslant \frac{1}{\alpha}\left\|\mathbf{f}\left(\mathbf{u}_{0}\right)\right\|_{V^{\prime}}
$$

By the triangular inequality, estimate (9) holds with $M=\left\|\mathbf{u}_{0}\right\|_{\mathbf{H}^{1}(\mathcal{F})}+1 / \alpha\left\|\mathbf{f}\left(\mathbf{u}_{0}\right)\right\|_{V^{\prime}}$. This concludes the proof of Theorem 1.

Proof of Theorem 2. Let us fix $\mathbf{u}_{\infty} \in \mathbb{R}^{N}$.
To prove $(v)$, we observe that for $\mathbf{v}, \widetilde{\mathbf{v}}, \mathbf{v}_{1}, \mathbf{v}_{2} \in V$,

$$
a\left(\mathbf{v}, \mathbf{v}_{1}, \mathbf{v}_{2}\right)-a\left(\widetilde{\mathbf{v}}, \mathbf{v}_{1}, \mathbf{v}_{2}\right)=a_{1}\left(\mathbf{v}-\widetilde{\mathbf{v}}, \mathbf{v}_{1}, \mathbf{v}_{2}\right) .
$$

Introducing

$$
\begin{equation*}
\mathcal{N}=\sup _{\mathbf{v}, \mathbf{v}_{1}, \mathbf{v}_{2} \in V \backslash\{0\}} \frac{a_{1}\left(\mathbf{v}, \mathbf{v}_{1}, \mathbf{v}_{2}\right)}{\|\mathbf{v}\|_{\mathbf{H}^{1}(\mathcal{F})}\left\|\mathbf{v}_{1}\right\|_{\mathbf{H}^{1}(\mathcal{F})}\left\|\mathbf{v}_{2}\right\|_{\mathbf{H}^{1}(\mathcal{F})}}, \tag{34}
\end{equation*}
$$

we obtain that (27) holds with $L(\eta)=\mathcal{N}$ for all $\eta>0$.
We define

$$
\nu_{0}=\inf \left\{C_{2}\left[\rho\left(\mathbf{u}_{0}\right)+\left(\mathcal{N}\left\|\mathbf{f}\left(\mathbf{u}_{0}\right)\right\|_{V^{\prime}}\right)^{1 / 2}\right], \mathbf{u}_{0} \in K \text { satisfying }(20)\right\} .
$$

Estimate (10) results from the existence of a constant $C_{3}>0$, such that for every $\phi \in V$,

$$
\left|\left\langle\mathbf{f}\left(\mathbf{u}_{0}\right), \phi\right\rangle_{V^{\prime} \times V}\right| \leqslant C_{3}\left\|\mathbf{u}_{0}\right\|_{\mathbf{H}^{1}(\mathcal{F})}\left(1+\left\|\mathbf{u}_{0}\right\|_{\mathbf{H}^{1}(\mathcal{F})}\right)\|\phi\|_{\mathbf{H}^{1}(\mathcal{F})} .
$$

As a result, by (20), there exists $C_{4}>0$ such that

$$
\left\|\mathbf{f}\left(\mathbf{u}_{0}\right)\right\|_{V^{\prime}} \leqslant C_{4}\left|\mathbf{u}_{\infty}\right|\left(1+\left|\mathbf{u}_{\infty}\right|\right),
$$

and consequently, inequality (35) implies that there exists a constant $C_{5}>0$ such that

$$
\nu_{0}^{2} \leqslant C_{5}\left(\left|\mathbf{u}_{\infty}\right|+\left|\mathbf{u}_{\infty}\right|^{2}\right) .
$$

This proves (10).
Let $\nu>\nu_{0}$. By definition of the infimum, there exists $\mathbf{u}_{0} \in K$ satisfying (20) and such that

$$
\begin{equation*}
\nu_{0} \leqslant C_{2}\left[\rho\left(\mathbf{u}_{0}\right)+\left(\mathcal{N}\left\|\mathbf{f}\left(\mathbf{u}_{0}\right)\right\|_{V^{\prime}}\right)^{1 / 2}\right]<\nu \tag{35}
\end{equation*}
$$

We fix $\alpha=\nu / C_{2}-\rho\left(\mathbf{u}_{0}\right)$. By (35), $\alpha>0$ and by the same argument as above, $\alpha$ satisfies (26). Now let $\mathbf{u} \in K$ be a weak solution to System (1)-(5) and let $\widetilde{\mathbf{u}}=\mathbf{u}-\mathbf{u}_{0}$. Condition (28) reads

$$
\mathcal{N}\left\|\mathbf{f}\left(\mathbf{u}_{0}\right)\right\|_{V^{\prime}}<\alpha^{2}
$$

which is equivalent to

$$
\begin{equation*}
\left(\mathcal{N}\left\|\mathbf{f}\left(\mathbf{u}_{0}\right)\right\|_{V^{\prime}}\right)^{1 / 2}<\nu / C_{2}-\rho\left(\mathbf{u}_{0}\right) \tag{36}
\end{equation*}
$$

By (35), last estimate holds. By Theorem 6, $\widetilde{\mathbf{u}} \in V$ is the unique solution to (23). Using (10) and (36), we obtain the following estimate:

$$
\begin{equation*}
\|\widetilde{\mathbf{u}}\|_{\mathbf{H}^{1}(\mathcal{F})}<\frac{1}{\mathcal{N}}\left(\frac{\nu}{C_{2}}-\rho\left(\mathbf{u}_{0}\right)\right) \tag{37}
\end{equation*}
$$

Consequently, be (20) and the triangular inequality,

$$
\|\mathbf{u}\|_{\mathbf{H}^{1}(\mathcal{F})} \leqslant C_{0}\left|\mathbf{u}_{\infty}\right|+\frac{1}{\mathcal{N} C_{2}} \nu .
$$

This proves (11).
By taking the infimum over all admissible functions $\mathbf{u}_{0}$, for $\nu>\nu_{0}$, we obtain the uniqueness of the function $\widetilde{\mathbf{u}} \in K$ satisfying (21). Therefore, System (1)-(5) admits one unique weak solution $\mathbf{u}_{0}+\widetilde{\mathbf{u}} \in \mathbf{H}^{1}(\mathcal{F})$.

## 4 Proof of Theorem 3

In this section, we assume that $\mathcal{V}(x)$ is equal to the tangent hyperplane to $\partial \mathcal{S}$ on q.e. $x$ on the boundary, $A \equiv \operatorname{Id}$ and $\mu$ is absolutely continuous with respect to the surface measure on $\partial \mathcal{S}$ with a density $\beta \in L^{2}(\partial \mathcal{S})$.

Let $\mathbf{u}_{\infty} \in \mathbb{R}^{N}$. Let $\nu_{0}>0$ as in Theorem 2, and $\nu>\nu_{0}$. Let $\mathbf{u}_{0} \in K$ satisfying (20) and such that

$$
\begin{equation*}
\nu_{0} \leqslant C_{2}\left[\rho\left(\mathbf{u}_{0}\right)+\left(\mathcal{N}\left\|\mathbf{f}\left(\mathbf{u}_{0}\right)\right\|_{V_{\tau}^{\prime}}\right)^{1 / 2}\right]<\nu \tag{38}
\end{equation*}
$$

We use the mappings introduced in Section 3. We define $\mathbf{g}=\mathbf{g}\left(\mathbf{u}_{0}\right) \in V_{\tau}^{\prime}$ by

$$
\langle\mathbf{g}, \phi\rangle_{V_{\tau}^{\prime} \times V_{\tau}}=-a_{0}\left(\mathbf{u}_{0}, \phi\right)-a_{1}\left(\mathbf{u}_{0}, \mathbf{u}_{0}, \phi\right) \quad \forall \phi \in V_{\tau}
$$

The problem we consider is equivalent to the following: find $\mathbf{u}_{0} \in K_{\tau}$ and $\widetilde{\mathbf{u}} \in V_{\tau}$ such that

$$
\begin{equation*}
a(\widetilde{\mathbf{u}}, \widetilde{\mathbf{u}}, \phi)=\left\langle\mathbf{g}\left(\mathbf{u}_{0}\right), \phi\right\rangle_{V_{\tau}^{\prime} \times V_{\tau}} \quad \forall \phi \in V_{\tau} \tag{39}
\end{equation*}
$$

Proof of Theorem 3. The proof is based on the implicit function theorem. We will prove that the mapping

$$
\beta \in \mathcal{O} \rightarrow \widetilde{\mathbf{u}} \in V_{\tau}
$$

where $\widetilde{\mathbf{u}}$ satisfies (21), is differentiable. We introduce the open subset of $V_{\tau}$ :

$$
V_{0}=\left\{\xi \in V_{\tau} \left\lvert\,\|\xi\|_{\mathbf{H}^{1}(\mathcal{F})}<\frac{1}{\mathcal{N}}\left(\frac{\nu}{C_{2}}-\rho\left(\mathbf{u}_{0}\right)\right)\right.\right\}
$$

and define the mapping

$$
F: \mathcal{O} \times V_{0} \rightarrow V_{\tau}^{\prime}
$$

by

$$
\begin{align*}
\forall \phi \in V_{\tau} \quad\langle F(\beta, \mathbf{v}), \phi\rangle_{V_{\tau}^{\prime} \times V_{\tau}}= & a_{0}(\mathbf{v}, \phi)+a_{1}\left(\mathbf{u}_{0}, \mathbf{v}, \phi\right)+a_{1}\left(\mathbf{v}, \mathbf{u}_{0}, \phi\right)+a_{1}(\mathbf{v}, \mathbf{v}, \phi) \\
& +a_{0}\left(\mathbf{u}_{0}, \phi\right)+a_{1}\left(\mathbf{u}_{0}, \mathbf{u}_{0}, \phi\right) . \tag{40}
\end{align*}
$$

The mapping $F$ is well defined and $F$ is of class $\mathcal{C}^{\infty}$. By the proof of Theorems 1 and 2, for every $\beta \in \mathcal{O}$ there exists one unique $\widetilde{\mathbf{u}} \in V_{0}$ satisfying (21). We need to prove that for every $(\beta, \mathbf{v}) \in \mathcal{O} \times V_{0}, D_{\mathbf{v}} F(\beta, \mathbf{v}) \in \mathcal{L}\left(V_{\tau}, V_{\tau}^{\prime}\right)$ is invertible with continuous inverse.

For every $(\beta, \mathbf{v}) \in \mathcal{O} \times V_{0}, D_{\mathbf{v}} F(\beta, \mathbf{v}) \in \mathcal{L}\left(V_{\tau}, V_{\tau}^{\prime}\right)$ is defined by

$$
\begin{aligned}
& \forall \xi, \phi \in V_{\tau} \quad\left\langle D_{\mathbf{v}} F(\beta, \mathbf{v}) \xi, \phi\right\rangle_{V_{\tau}^{\prime} \times V_{\tau}}=a_{0}(\xi, \phi)+a_{1}\left(\mathbf{u}_{0}, \xi, \phi\right)+a_{1}\left(\xi, \mathbf{u}_{0}, \phi\right) \\
&+a_{1}(\xi, \mathbf{v}, \phi)+a_{1}(\mathbf{v}, \xi, \phi)
\end{aligned}
$$

We need to prove that for every $(\beta, \mathbf{v}) \in \mathcal{O} \times V_{0}$, there exists a constant $C>0$ such that for every $\mathbf{h} \in V_{\tau}^{\prime}$, there exists one unique $\xi \in V_{\tau}$ such that

$$
\begin{equation*}
D_{\mathbf{v}} F(\beta, \mathbf{v}) \xi=\mathbf{h} \tag{41}
\end{equation*}
$$

with the estimate

$$
\begin{equation*}
\|\xi\|_{\mathbf{H}^{1}(\mathcal{F})} \leqslant C\|\mathbf{h}\|_{V_{\tau}^{\prime}} . \tag{42}
\end{equation*}
$$

Fix $(\beta, \mathbf{v}) \in \mathcal{O} \times V_{0}$. We will prove that the bilinear form associated to $D_{\mathbf{v}} F(\beta, \mathbf{v})$ satisfies the assumptions of Lax-Milgram theorem. The continuity results from the continuity of $a_{0}$ and $a_{1}$, and the coercivity is a consequence of the definition of $V_{0}$. Indeed, by computation (31), for all $\phi \in V_{\tau}$,

$$
\left\langle D_{\mathbf{v}} F(\beta, \mathbf{v}) \phi, \phi\right\rangle_{V_{\tau}^{\prime} \times V_{\tau}}=a_{0}(\phi, \phi)+a_{1}\left(\phi, \mathbf{u}_{0}, \phi\right)+a_{1}(\phi, \mathbf{v}, \phi)
$$

Consequently, using property (30) and definitions (32) and (34), we obtain for all $\phi \in V_{\tau}$,

$$
\left\langle D_{\mathbf{v}} F(\beta, \mathbf{v}) \phi, \phi\right\rangle_{V_{\tau}^{\prime} \times V_{\tau}} \geqslant\left(\frac{\nu}{C_{2}}-\rho\left(\mathbf{u}_{0}\right)-\mathcal{N}\|v\|_{\mathbf{H}^{1}(\mathcal{F})}\right)\|\phi\|_{\mathbf{H}^{1}(\mathcal{F})}^{2} .
$$

Since $\mathbf{v} \in V_{0}$,

$$
\frac{\nu}{C_{2}}-\rho\left(\mathbf{u}_{0}\right)-\mathcal{N}\|v\|_{\mathbf{H}^{1}(\mathcal{F})}>0 .
$$

This proves that the mapping

$$
\beta \in \mathcal{O} \rightarrow \widetilde{\mathbf{u}} \in V_{\tau}
$$

where $\widetilde{\mathbf{u}}$ satisfies (39), is differentiable. Consequently, since for all $\beta \in \mathcal{O}, R(\beta)=\mathbf{u}_{0}+\widetilde{\mathbf{u}}, R$ is differentiable on $\mathcal{O}$. Moreover, for every $\beta \in \mathcal{O}, D R(\beta) \in \mathcal{L}\left(L^{2}(\partial \mathcal{S}), \mathbf{H}\right)$ is defined by

$$
D R(\beta)=-\left[D_{v} F\left(\beta, R(\beta)-\mathbf{u}_{0}\right)\right]^{-1} D_{\beta} F\left(\beta, R(\beta)-\mathbf{u}_{0}\right)
$$

Since for every $(\beta, \mathbf{v}) \in \mathcal{O} \times V_{0}, D_{\beta} F(\beta, \mathbf{v}) \in \mathcal{L}\left(L^{2}(\partial \mathcal{S}), V_{\tau}^{\prime}\right)$ is defined by

$$
\forall h \in L^{2}(\partial \mathcal{S}), \forall \phi \in V_{\tau} \quad\left\langle D_{\beta} F(\beta, \mathbf{v}), h\right\rangle=\int_{\partial \mathcal{S}} h\left(\mathbf{v}+\mathbf{u}_{0}\right) \cdot \phi \mathrm{d} \mathcal{H}^{N-1}
$$

defining $\mathbf{w}=D R(\beta) h$ for $h \in L^{2}(\partial \mathcal{S})$, we obtain that $\mathbf{w} \in V_{\tau}$ satisfies the following identity:

$$
\begin{aligned}
\forall \phi \in V_{\tau} \quad & a_{0}(\mathbf{w}, \phi)+a_{1}\left(\mathbf{u}_{0}, \mathbf{w}, \phi\right)+a_{1}\left(\mathbf{w}, \mathbf{u}_{0}, \phi\right)+a_{1}\left(\mathbf{w}, R(\beta)-\mathbf{u}_{0}, \phi\right) \\
& +a_{1}\left(R(\beta)-\mathbf{u}_{0}, \mathbf{w}, \phi\right)=-\int_{\partial \mathcal{S}} h R(\beta) \cdot \phi \mathrm{d} \mathcal{H}^{N-1} .
\end{aligned}
$$

Using the trilinearity of $a_{1}$, we obtain (14).
Proof of Corollary 4. Noticing that for every $\phi \in V_{\tau}$,

$$
\left[(\nabla \mathbf{u})^{T} \psi\right] \cdot \phi=[(\phi \cdot \nabla) \mathbf{u}] \cdot \psi
$$

and

$$
-\int_{\mathcal{F}}[(\mathbf{u} \cdot \nabla) \psi] \cdot \phi \mathrm{d} x=\int_{\mathcal{F}}[(\mathbf{u} \cdot \nabla) \phi] \cdot \psi \mathrm{d} x
$$

(17) can be written in the following form:

$$
\begin{align*}
\forall \phi \in V_{\tau} \quad & 2 \nu \int_{\mathcal{F}} \mathbf{D}(\psi): \mathbf{D}(\phi) \mathrm{d} x+\int_{\mathcal{F}} \psi \cdot[(\phi \cdot \nabla) \mathbf{u}+(\mathbf{u} \cdot \nabla) \phi] \mathrm{d} x+\int_{\partial \mathcal{S}} \beta \psi \cdot \phi \mathrm{d} \mathcal{H}^{N-1} \\
& =2 \int_{\mathcal{F}}[(\mathbf{u} \cdot \nabla) \mathbf{u}] \cdot \phi \mathrm{d} x . \tag{43}
\end{align*}
$$

Let us fix $\beta \in \mathcal{O}$ and $h \in L^{2}(\partial \mathcal{S})$. We note $\mathbf{u}=R(\beta)$ and $\mathbf{w}=D R(\beta) h$. Taking $\phi=\mathbf{w}$ in (12), we express $D T(\beta) h$ as follows:

$$
\begin{equation*}
D T(\beta) h=-2 \int_{\mathcal{F}}[(\mathbf{u} \cdot \nabla) \mathbf{u}] \cdot \mathbf{w} \mathrm{d} x+\int_{\partial \mathcal{S}} h|\mathbf{u}|^{2} \mathrm{~d} \mathcal{H}^{N-1} \tag{44}
\end{equation*}
$$

Besides, taking $\phi=\psi$ in (14) and $\phi=\mathbf{w}$ in (43), we obtain

$$
\begin{aligned}
& 2 \nu \int_{\mathcal{F}} \mathbf{D}(\mathbf{w}): \mathbf{D}(\psi) \mathrm{d} x+\int_{\mathcal{F}}[(\mathbf{w} \cdot \nabla) \mathbf{u}+(\mathbf{u} \cdot \nabla) \mathbf{w}] \cdot \psi \mathrm{d} x+\int_{\partial \mathcal{S}} \beta \mathbf{w} \cdot \psi \mathrm{d} \mathcal{H}^{N-1} \\
& \quad=-\int_{\partial \mathcal{S}} h \mathbf{u} \cdot \psi \mathrm{~d} \mathcal{H}^{N-1}, \\
& 2 \nu \int_{\mathcal{F}} \mathbf{D}(\psi): \mathbf{D}(\mathbf{w}) \mathrm{d} x+\int_{\mathcal{F}} \psi \cdot[(\mathbf{w} \cdot \nabla) \mathbf{u}+(\mathbf{u} \cdot \nabla) \mathbf{w}] \mathrm{d} x+\int_{\partial \mathcal{S}} \beta \psi \cdot \mathbf{w} \mathrm{d} \mathcal{H}^{N-1} \\
& \quad=2 \int_{\mathcal{F}}[(\mathbf{u} \cdot \nabla) \mathbf{u}] \cdot \mathbf{w} \mathrm{d} x .
\end{aligned}
$$

Substracting these equalities, we obtain

$$
-2 \int_{\mathcal{F}}[(\mathbf{u} \cdot \nabla) \mathbf{u}] \cdot \mathbf{w} \mathrm{d} x=\int_{\partial \mathcal{S}} h \mathbf{u} \cdot \psi \mathrm{~d} \mathcal{H}^{N-1}
$$

and finally, by (44),

$$
D T(\beta) h=\int_{\partial \mathcal{S}}[(\mathbf{u}+\psi) \cdot \mathbf{u}] h \mathrm{~d} \mathcal{H}^{N-1} .
$$

This proves (15).
Remark 9 The coercivity of the bilinear form associated to the weak formulation (43) results from the following estimate (see the proof of Theorems 1 and 2):

$$
\begin{aligned}
\rho(\mathbf{u}) & \leqslant \rho\left(\mathbf{u}_{0}\right)+\rho(\widetilde{\mathbf{u}}) \\
& \leqslant \rho\left(\mathbf{u}_{0}\right)+\mathcal{N}\|\widetilde{\mathbf{u}}\|_{\mathbf{H}^{1}(\mathcal{F})} \\
& <\frac{\nu}{C_{2}} .
\end{aligned}
$$

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