

Tangential Navier–Stokes equations on evolving surfaces: Analysis and simulations

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The paper considers a system of equations that models a lateral flow of a Boussinesq–Scriven fluid on a passively evolving surface embedded in \mathbb{R}^3 . For the resulting Navier–Stokes type system, posed on a smooth closed time-dependent surface, we introduce a weak formulation in terms of functional spaces on a space-time manifold defined by the surface evolution. The weak formulation is shown to be well-posed for any finite final time and without smallness conditions on data. We further extend an unfitted finite element method, known as TraceFEM, to compute solutions to the fluid system. Convergence of the method is demonstrated numerically. In another series of experiments we visualize lateral flows induced by smooth deformations of a material surface.

Keywords: Navier-Stokes equations; well-posedness; finite element method.

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1. Introduction

There is extensive literature on analysis and numerical simulation of the incompressible Navier–Stokes equations, a fundamental model of fluid mechanics. While

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the overwhelming majority of papers in this field treats these equations in Euclidean domains, there also is literature on analysis of the incompressible Navier–Stokes equations on surfaces, or more general on Riemannian manifolds. Building on a fundamental observation made by Arnold² that relates equations of incompressible fluid to finding geodesics on the group of all volume preserving diffeomorphisms, local existence and uniqueness results for Navier–Stokes equations on compact oriented Riemannian manifolds were proved in the seminal paper (Ref. 11). This work has been followed by many other studies, cf. Refs. 44 and 42 and the overview in Ref. 8. Very recent activity in the field includes the work,^{33, 40} in which local-in-time-well-posedness in the framework of maximal regularity is established. All these papers restrict to *stationary* surfaces or manifolds.

In recent years there has been a growing interest in fluid equations on evolving surfaces.^{16, 20, 24, 38, 47} motivated in particular by applications to modeling of biological membranes, e.g. see Refs. 39, 34, 4 and 45. In Ref. 6 one finds an overview and comparison of different modeling approaches for evolving viscous fluid lavers that result in the surface Navier–Stokes equations. We are not aware of any literature presenting well-posedness analysis of this system on evolving surfaces. Furthermore, only very few papers address numerical treatment of such equations. In Refs. 36 and 37 computational results are presented, based on a surface vorticitystream function formulation of the Navier–Stokes equations. The surface motion is prescribed and the evolving SFEM of Dziuk–Elliott^{9, 10} is applied to the partial differential equations for the scalar vorticity and stream function unknowns. The authors of Ref. 26 consider another discretization approach that is based on the techniques developed in Ref. 38. These papers focus on modeling and illustration of certain interesting flow phenomena but not on the performance of the numerical methods. Several recent papers^{5, 14, 17, 19} present error analysis of finite element discretization methods for vector-valued PDEs on stationary surfaces. We are not aware of any paper with a systematic numerical study or an error analysis of a discretization method for vector-valued PDEs on evolving surfaces. We conclude that in the field of incompressible Navier–Stokes equations on time-dependent surfaces basic problems related to well-posedness of the systems, development and analysis of numerical methods remain open. This paper addresses two of these problems: well-posedness and numerical method development.

It is shown in Ref. 6 that several different modeling approaches all yield the same *tangential* surface Navier–Stokes equations (TSNSE). These equations govern the evolution of tangential velocity and surface pressure if the normal velocity of the surface is prescribed. The main topic of this paper is the analysis of a variational formulation of the TSNSE. In particular, a well-posedness result for this formulation is proved. To the best of our knowledge, this is the first well-posedness result for evolving surface Navier–Stokes equations. The paper also touches on the development of a new discretization method for the TSNSE. This method combines an implicit time stepping scheme with a TraceFEM^{28, 29} for discretization in space. We explain this method, validate its optimal second-order convergence for a test

problem with a known solution and apply it to the simulation of a lateral flow induced by deformations of a sphere. Error analysis of this method is not addressed in this paper and left for future research.

The remainder of the paper is organized as follows. In Sec. 2, we recall the surface Navier–Stokes equations known from the literature. In particular, the TSNSEs are described. Appropriate function spaces for a variational formulation of the TSNSE are introduced in Sec. 3. Relevant properties of these spaces are derived. The main results of this paper are given in Sec. 4. We introduce and analyze two variational formulations of the TSNE: The first one is for the tangential velocity only, which is solenoidal by construction of the solution space. Then we introduce the pressure and study a mixed variational problem. For both formulations well-posedness results are derived. In Sec. 5, we explain a discretization method. Finally, Sec. 6 collects and discusses results of numerical experiments.

2. Surface Navier–Stokes Equations

We first introduce necessary notations of surface quantities and tangential differential operators. For a closed smooth surface Γ embedded in \mathbb{R}^3 , the outward pointing normal vector is denoted by \mathbf{n} , and the normal projector on the tangential space at $\mathbf{x} \in \Gamma$ is $\mathbf{P} = \mathbf{P}(\mathbf{x}) = \mathbf{I} - \mathbf{nn}^T$. Let $\mathbf{H} = \nabla_{\Gamma} \mathbf{n} \in \mathbb{R}^{3\times 3}$ be the Weingarten mapping (second fundamental form) and $\kappa := \operatorname{tr}(\mathbf{H})$ twice the mean curvature. For a scalar function $p: \Gamma \to \mathbb{R}$ or a vector field $\mathbf{u}: \Gamma \to \mathbb{R}^3$ their smooth extensions to a neighborhood $\mathcal{O}(\Gamma)$ of Γ are denoted by p^e and \mathbf{u}^e , respectively. Surface gradients and covariant derivatives on Γ can be defined through derivatives in \mathbb{R}^3 as $\nabla_{\Gamma} p = \mathbf{P} \nabla p^e$, $D_{\Gamma} \mathbf{u} \coloneqq \nabla \mathbf{u}^e \mathbf{P}$, and $\nabla_{\Gamma} \mathbf{u} \coloneqq \mathbf{P} \nabla \mathbf{u}^e \mathbf{P}$. These definitions are independent of a particular smooth extension of p and \mathbf{u} off Γ . The surface rate-of-strain tensor¹³ is given by $E_s(\mathbf{u}) \coloneqq \frac{1}{2} (\nabla_{\Gamma} \mathbf{u} + \nabla_{\Gamma} \mathbf{u}^T)$, the surface divergence and $\operatorname{curl}_{\Gamma}$ operators for a vector field $\mathbf{u}: \Gamma \to \mathbb{R}^3$ are $\operatorname{div}_{\Gamma} \mathbf{u} \coloneqq \operatorname{tr}(\nabla_{\Gamma} \mathbf{u})$ and $\operatorname{curl}_{\Gamma} \mathbf{u} \coloneqq (\nabla_{\Gamma} \times \mathbf{u}) \cdot \mathbf{n}$. For a tensor field $\mathbf{A} = [\mathbf{a}_1, \mathbf{a}_2, \mathbf{a}_3]: \Gamma \to \mathbb{R}^{3\times 3}$, $\operatorname{div}_{\Gamma} \mathbf{A}$ is defined row-wise and $D_{\Gamma} \mathbf{A}$ is a third-order tensor such that $(D_{\Gamma} \mathbf{A})_{i,j,k} = (D_{\Gamma} \mathbf{a}_j)_{i,k}$.

We now let $\Gamma(t)$ be a material surface embedded in \mathbb{R}^3 as defined in Refs. 13 and 25, with a density distribution $\rho(t, \mathbf{x})$. By $\mathbf{u}(t, \mathbf{x}), \mathbf{x} \in \Gamma(t)$, we denote a velocity field of the density flow on Γ , i.e. $\mathbf{u}(t, \mathbf{x})$ is the velocity of a material point $\mathbf{x} \in \Gamma(t)$. The derivative f of a surface quantity f along the corresponding trajectories of material points is called the material derivative. Assuming the surface evolution is such that the space-time manifold

$$\mathcal{S} = \bigcup_{t \in [0,T]} \{t\} \times \Gamma(t) \subset \mathbb{R}^4$$

is smooth, the material derivative can be defined as

$$\dot{f} = \frac{\partial f^e}{\partial t} + (\mathbf{u} \cdot \nabla) f^e \quad \text{on } \mathcal{S},$$
(2.1)

where f^e is a smooth extension of $f : S \to \mathbb{R}$ into a spatial neighborhood of S. Note that \dot{f} is a tangential derivative for S, and hence it depends only on the surface values of f on S. For a vector field \mathbf{v} on S, one defines $\dot{\mathbf{v}}$ componentwise.

The conservation of mass and linear momentum for a thin material layer represented by $\Gamma(t)$ together with the Boussinesq–Scriven constitutive relation for the surface stress tensor and an inextensibility condition leads to the *surface Navier–Stokes equations*:

$$\begin{cases} \rho \dot{\mathbf{u}} = -\nabla_{\Gamma} \pi + 2\mu \operatorname{div}_{\Gamma}(E_{s}(\mathbf{u})) + \mathbf{b} + \pi \kappa \mathbf{n}, \\ \operatorname{div}_{\Gamma} \mathbf{u} = 0, & \text{on } \Gamma(t), \\ \dot{\rho} = 0, \end{cases}$$
(2.2)

where π is the surface fluid pressure and μ stands for the viscosity coefficient. Equations (2.2) model the evolution of an inextensible viscous fluidic material surface with acting area force **b**, cf. Refs. 20 and 16 for derivations of this model and Ref. 6 for a literature overview and alternative forms of this system. The pure geometrical evolution of $\Gamma(t)$ is defined by its normal velocity $V_{\Gamma} = V_{\Gamma}(t, \mathbf{x})$ that is given by the normal component of the material velocity,

$$V_{\Gamma} = \mathbf{u} \cdot \mathbf{n} \quad \text{on } \Gamma(t). \tag{2.3}$$

If **b** is given or defined through other unknowns, then (2.2) and (2.3) form a closed system of six equations for six unknowns \mathbf{u} , π , ρ , and V_{Γ} , subject to suitable initial conditions.

2.1. Tangential surface Navier-Stokes equations

We now introduce a major simplification by assuming that the *geometric* evolution of Γ is known. We make this more precise below and derive equations governing the unknown lateral motions of the surface fluid. To this end, consider a smooth velocity field $\mathbf{w} = \mathbf{w}(t, \mathbf{x})$ in $[0, T] \times \mathbb{R}^3$ that passively advects the embedded surface $\Gamma(t)$ given by

$$\Gamma(t) = \{ \mathbf{y} \in \mathbb{R}^3 \, | \, \mathbf{y} = \mathbf{x}(t, \mathbf{z}), \mathbf{z} \in \Gamma_0 \},$$
(2.4)

where the trajectories $\mathbf{x}(t, \mathbf{z})$ are the unique solutions of the Cauchy problem

$$\begin{cases} \mathbf{x}(0, \mathbf{z}) = \mathbf{z}, \\ \frac{d}{dt} \mathbf{x}(t, \mathbf{z}) = \mathbf{w}(t, \mathbf{x}(t, \mathbf{z})), \end{cases}$$
(2.5)

for all \mathbf{z} on an initial smooth connected surface $\Gamma_0 = \Gamma(0)$ embedded in \mathbb{R}^3 . We now assume that the normal material motion of Γ is completely determined by the ambient flow \mathbf{w} and the lateral material motion is free, i.e. for the *given* \mathbf{w} the relation

$$\mathbf{u} \cdot \mathbf{n} = \mathbf{w} \cdot \mathbf{n} \quad \text{on } \Gamma(t) \tag{2.6}$$

holds for the normal component^a $u_N = \mathbf{u} \cdot \mathbf{n}$, while the tangential component \mathbf{u}_T of the surface fluid flow is unknown and depends on \mathbf{w} only implicitly through the variation of $\Gamma(t)$ and conservation laws represented by Eq. (2.2). The resulting system can be seen as an idealized model for the motion of a fluid layer embedded in bulk fluid, where one neglects friction forces between the surface and the bulk as well as any effect of the layer on the bulk flow. In such a physical setting, (2.6) means non-penetration of the bulk fluid through the material layer.

Material trajectories of points on the surface are defined by the flow field \mathbf{u} , rather than \mathbf{w} . We are also interested in a derivative determined by the variation of a quantity along the so-called normal trajectories defined below.

Definition 2.1. Let $\Phi_t^n : \Gamma_0 \to \Gamma(t)$, be the flow map of the pure geometric (normal) evolution of the surface, i.e. for $\mathbf{z} \in \Gamma_0$, the normal trajectory $\mathbf{x}^n(t, \mathbf{z}) = \Phi_t^n(\mathbf{z})$ solves

$$\begin{cases} \mathbf{x}^{n}(0, \mathbf{z}) = \mathbf{z}, \\ \frac{d}{dt} \mathbf{x}^{n}(t, \mathbf{z}) = \mathbf{w}_{N}(t, \mathbf{x}^{n}(\mathbf{z}, t)). \end{cases}$$
(2.7)

Equation (2.7) defines a bijection between Γ_0 and $\Gamma(t)$ for every $t \in [0,T]$ with inverse mapping Φ^n_{-t} . The Lagrangian derivative for the flow map Φ^n_t is denoted by ∂° :

$$\partial^{\circ} \mathbf{v}(t, \mathbf{x}) = \frac{d}{dt} \mathbf{v}(t, \Phi_t^n(\mathbf{z})), \quad \mathbf{x} = \Phi_t^n(\mathbf{z}).$$
(2.8)

We call $\partial^{\circ} \mathbf{v}$ the normal time derivative of \mathbf{v} .

It is clear from (2.8) that this normal time derivative is an intrinsic surface quantity. Similar to the material derivative in (2.1), it can be expressed in terms of bulk derivatives if one assumes a smooth extension of \mathbf{v} from \mathcal{S} to its neighborhood:

$$\partial^{\circ} \mathbf{v}(t, \mathbf{x}) = \frac{d}{dt} \mathbf{v}(t, \mathbf{x}^{n}(t, \mathbf{z})) = \left(\frac{\partial \mathbf{v}^{e}}{\partial t} + (\mathbf{w}_{N} \cdot \nabla) \mathbf{v}^{e}\right)(t, \mathbf{x})$$
(2.9)

for $(t, \mathbf{x}) \in S$. Comparing the material and normal time derivatives of a flow field **v** on the surface we find the equality

$$\mathbf{P}\mathbf{\dot{v}} = \mathbf{P}\partial^{\circ}\mathbf{v} + (\nabla_{\Gamma}\mathbf{v})\mathbf{u}_{T}.$$

With the splitting $\mathbf{v} = \mathbf{v}_T + \mathbf{v}_N$, we get

$$\mathbf{P}\dot{\mathbf{v}} = \mathbf{P}\partial^{\circ}\mathbf{v}_{T} + \mathbf{P}\partial^{\circ}\mathbf{v}_{N} + (\nabla_{\Gamma}\mathbf{v}_{T})\mathbf{u}_{T} + (\nabla_{\Gamma}\mathbf{v}_{N})\mathbf{u}_{T}.$$
 (2.10)

Noting that $\mathbf{Pn} = 0$ and $\mathbf{P}\partial^{\circ}\mathbf{n} = -\nabla_{\Gamma}w_N$ (cf. (2.16) in Ref. 16), we rewrite $\mathbf{P}\partial^{\circ}\mathbf{v}_N$ as

$$\mathbf{P}\partial^{\circ}\mathbf{v}_{N}=\partial^{\circ}v_{N}\mathbf{P}\mathbf{n}+v_{N}\mathbf{P}\partial^{\circ}\mathbf{n}=v_{N}\mathbf{P}\partial^{\circ}\mathbf{n}=-v_{N}\nabla_{\Gamma}w_{N}\mathbf{x}$$

^aFor velocity fields $\mathbf{v} \in \mathbb{R}^3$ defined on $\Gamma(t)$ we use a splitting into tangential and normal components $\mathbf{v} = \mathbf{v}_T + \mathbf{v}_N = \mathbf{v}_T + v_N \mathbf{n}$, with $v_N = \mathbf{v} \cdot \mathbf{n}$.

We also have the relation $(\nabla_{\Gamma} \mathbf{v}_N) \mathbf{u}_T = v_N \mathbf{H} \mathbf{u}_T$. Using these results and letting $\mathbf{v} = \mathbf{u}$ in (2.10) one obtains

$$\mathbf{P}\dot{\mathbf{u}} = \mathbf{P}\partial^{\circ}\mathbf{u}_{T} + (\nabla_{\Gamma}\mathbf{u}_{T})\mathbf{u}_{T} + w_{N}\mathbf{H}\mathbf{u}_{T} - \frac{1}{2}\nabla_{\Gamma}w_{N}^{2}, \qquad (2.11)$$

where we also used $u_N = w_N$. To derive an equation for the unknown tangential velocity \mathbf{u}_T , we apply the projection \mathbf{P} to the first equation in (2.2). For $\mathbf{P}\dot{\mathbf{u}}$ we have the result (2.11). Note that the term $\frac{1}{2}\nabla_{\Gamma}w_N^2$ is known and can be treated as a source term. For a stationary surface $(w_N = 0)$ the normal time derivative is just the usual time derivative, $\mathbf{P}\partial^{\circ}\mathbf{u}_T = \frac{\partial \mathbf{u}_T}{\partial t}$. The term $(\nabla_{\Gamma}\mathbf{u}_T)\mathbf{u}_T$ is the analog of the quadratic term in the Navier–Stokes equations. Using $\operatorname{div}_{\Gamma}\mathbf{u}_N = u_N\kappa$ and $u_N = w_N$, the second equation in (2.2) yields $\operatorname{div}_{\Gamma}\mathbf{u}_T = -w_N\kappa$. We are not interested in variable density case and let $\rho = 1$. Summarizing, from the surface Navier–Stokes equations (2.2) we get the following reduced system for \mathbf{u}_T and π which we call the TSNSE:

$$\begin{cases} \mathbf{P}\partial^{\circ}\mathbf{u}_{T} + (\nabla_{\Gamma}\mathbf{u}_{T})\mathbf{u}_{T} + w_{N}\mathbf{H}\mathbf{u}_{T} - 2\mu\mathbf{P}\operatorname{div}_{\Gamma}E_{s}(\mathbf{u}_{T}) + \nabla_{\Gamma}\pi = \mathbf{f},\\ \operatorname{div}_{\Gamma}\mathbf{u}_{T} = g, \end{cases}$$
(2.12)

with right-hand sides known in terms of geometric quantities, w_N and the tangential component of the external area force **b**:

$$g = -w_N \kappa, \quad \mathbf{f} = \mathbf{b}_T + 2\mu \mathbf{P} \operatorname{div}_{\Gamma}(w_N \mathbf{H}) + \frac{1}{2} \nabla_{\Gamma} w_N^2.$$
 (2.13)

In the remainder of this paper, we study this TSNSE. Note that these equations have a structure similar to the standard incompressible Navier–Stokes equations in Euclidean domains. Important differences are that TSNSE is formulated on a space-time manifold that does not have an evident tensor product structure and, related to this, a normal time derivative $\mathbf{P}\partial^{\circ}$ instead of the standard time derivative is used and an additional term $w_N \mathbf{H} \mathbf{u}_T$ occurs. After some preliminary results in the following section, we introduce a well-posed weak formulation of the TSNSE in Sec. 4.

Remark 2.1. If one does *not* assume a given normal velocity $u_N = w_N$, an equation for u_N can be derived from (2.2), cf. Ref. 16. The surface Navier–Stokes equations (2.2) are then rewritten as a *coupled system* for \mathbf{u}_T , π and u_N , that consists of TSNSE (2.12) and the coupled equation

$$\dot{\boldsymbol{u}}_N = -2\mu(\operatorname{tr}(\mathbf{H}\nabla_{\Gamma}\mathbf{u}_T) + \boldsymbol{u}_N\operatorname{tr}(\mathbf{H}^2)) + \mathbf{u}_T\cdot\mathbf{H}\mathbf{u}_T - \mathbf{u}_T\cdot\nabla_{\Gamma}\boldsymbol{u}_N + \pi\kappa + \boldsymbol{b}_N.$$
(2.14)

A challenging problem, not addressed in this paper, is the well-posedness of the surface Navier–Stokes equations (2.2), i.e. of the coupled system (2.12)–(2.14). For studying this problem, results on well-posedness of only the TSNSE (2.12) may be useful.

3. Preliminaries

In this section, we introduce several function spaces and derive relevant properties of these spaces. We will use these spaces to formulate a well-posed weak formulation of the TSNSE (2.12). At this point, we make our assumptions on Γ_0 and its evolution more precise. We introduce the following smoothness assumptions:

$$\Gamma_0 \in C^3$$
 and $\mathbf{w} \in C^3([0,T] \times \mathbb{R}^3, \mathbb{R}^3), \quad \sup_{[0,T] \times \mathbb{R}^3} |\mathbf{w}| < +\infty.$ (3.1)

Then the ODE system (2.5) has a unique solution for any $\mathbf{z} \in \Gamma_0 \subset \mathbb{R}^3$, which defines a one-to-one mapping $\Gamma_0 \to \Gamma(t)$ for all $t \in [0, T]$ (Theorems II.1.1, V.3.1 and remark to Theorem V.2.1 in Ref. 15). Moreover, this mapping is $C^3(\mathcal{S}_0, \mathbb{R}^4)$ (Corollary V.4.1 in Ref. 15) with

$$\mathcal{S}_0 := [0, T] \times \Gamma_0.$$

Therefore, S is a C^3 manifold as the image of $S_0 \in C^3$ under a smooth mapping.

We need a globally C^2 -smooth extension of the spatial normal $\mathbf{n}(t, x)$, $(t, x) \in S$ that can be constructed as follows. Let ϕ_0 be the signed distance function to Γ_0 . On a tubular neighborhood U_{δ} of Γ_0 , with diameter $\delta > 0$ sufficiently small, we have $\phi_0 \in C^3(U_{\delta})$, cf. Lemma 2.8 in Ref. 10. We extend this function to be from $C^3(\mathbb{R}^3)$ and zero outside $U_{2\delta}$. Thus, we have $\phi_0 \in C^3(\mathbb{R}^3)$ and ϕ_0 is a signed distance function in a neighborhood of Γ_0 . Let Φ_t be the flow map for the velocity field \mathbf{w} . The mapping $(t, \mathbf{x}) \to \Phi_t(\mathbf{x})$ is $C^3([0, T] \times \mathbb{R}^3, \mathbb{R}^3)$ and $\nabla_x \Phi_t(\mathbf{x})$ is regular.¹⁵ Define the level set function $\phi(t, x) := \phi_0(\Phi_{-t}(\mathbf{x}))$ and the neighborhood $S^{\text{ex}} := \bigcup_{t \in [0,T]} \{t\} \times$ $\Phi_t(U_{\delta})$ of S. Then we have $\phi \in C^3([0,T] \times \mathbb{R}^3, \mathbb{R})$ and for $(t, \mathbf{x}) \in S^{\text{ex}}$ it holds $|\nabla \phi(t, x)| \ge c > 0$, and $\phi(t, x) = 0$ iff $(t, x) \in S$. Set $\hat{\mathbf{n}}(t, \mathbf{x}) := \nabla \phi(t, x)/|\nabla \phi(t, x)|$ for $(t, \mathbf{x}) \in S^{\text{ex}}$. Clearly $\hat{\mathbf{n}} = \mathbf{n}$ on S and $\hat{\mathbf{n}} \in C^2(S^{\text{ex}}, \mathbb{R}^3)$, and by a standard procedure we can extend it to $\hat{\mathbf{n}} \in C^2([0,T] \times \mathbb{R}^3, \mathbb{R}^3)$. To simplify the notation, this extension is denoted by \mathbf{n} . For such an extended vector field \mathbf{n} we have that $\mathbf{w}_N = (\mathbf{w} \cdot \mathbf{n})\mathbf{n} \in C^2([0,T] \times \mathbb{R}^3, \mathbb{R}^3)$ holds. Arguing in the same way as above, we conclude that for the normal flow mapping from Definition 2.1, we have

$$\Phi_{(\cdot)}^n \in C^2(\mathcal{S}_0, \mathcal{S}). \tag{3.2}$$

Note that $S_0 = \overline{S_0}$ and $S = \overline{S}$, i.e. S_0 and S are closed manifolds.

We need function spaces suitable for a weak formulation of the TSNSE. For this we make use of a general framework of evolving spaces presented in Ref. 1. In Sec. 3.2, we introduce specific evolving Hilbert spaces, based on a Piola pushforward mapping. Based on results from Ref. 1 several properties of these spaces are derived. In Sec. 3.4, an evolving space of functions for which suitable weak "material" derivatives exist is introduced. Here we deviate from Ref. 1 in the sense that this "material" derivative is not based on the pushforward map but on the normal time derivative defined above.

3.1. Surface Piola transform

To define evolving Hilbert spaces based on standard Bochner spaces, we need a suitable pushforward map. In the context of this paper, it is natural to use a surface Piola transform as pushforward map, since this transform conserves the solenoidal property of a tangential vector field.

To define a surface variant of the Piola transform based on the normal flow map $\Phi_t^n : \Gamma_0 \to \Gamma(t)$, we need some further notation. Below we always take $\mathbf{z} \in \Gamma_0$ and $\mathbf{x} := \Phi_t^n(\mathbf{z}) \in \Gamma(t)$. Since for each $t \in [0,T]$ the map $\Phi_t^n : \Gamma_0 \to \Gamma(t)$ is a C^2 -diffeomorphism, the differential $D\Phi_t^n(\mathbf{z}) : (T\Gamma_0)_{\mathbf{z}} \to T\Gamma(t)_{\mathbf{x}}$, is invertible. Define $J = J(t, \mathbf{z}) := \det D\Phi_t^n(\mathbf{z}), J^{-1} = J^{-1}(t, \mathbf{x}) = \det D\Phi_{-t}^n(\mathbf{x}) = J(t, \mathbf{z})^{-1}$. Denote by $\mathbf{D} = \mathbf{D}(t, \mathbf{z})$ and $\mathbf{D}^{-1} = \mathbf{D}^{-1}(t, \mathbf{x})$ the matrices of linear mappings given by $D\Phi_t^n(\mathbf{z})\mathbf{P}(\mathbf{z}) : \mathbb{R}^3 \to \mathbb{R}^3$ and $D\Phi_{-t}^n(\mathbf{x})\mathbf{P}(\mathbf{x}) = [D\Phi_t^n(\mathbf{z})]^{-1}\mathbf{P}(\mathbf{x}) : \mathbb{R}^3 \to \mathbb{R}^3$, respectively. Note that $\mathbf{D}^{-1}\mathbf{D} = \mathbf{P}(\mathbf{z})$ and $\mathbf{D}\mathbf{D}^{-1} = \mathbf{P}(\mathbf{x})$ hold. For these mappings the following useful identities hold:

$$D_{\Gamma} \left(\mathbf{v} \circ \Phi_{-t}^{n} \right) = (D_{\Gamma_{0}} \mathbf{v}) \mathbf{D}^{-1} \quad \text{for } \mathbf{v} \in C^{1}(\Gamma_{0})^{3},$$

$$(D_{\Gamma} \mathbf{v}) \mathbf{D} = D_{\Gamma_{0}} (\mathbf{v} \circ \Phi_{t}^{n}) \quad \text{for } \mathbf{v} \in C^{1}(\Gamma(t))^{3}.$$
(3.3)

We need the Piola transform for arbitrary, not necessarily tangential vectors. For this it is convenient to define an invertible operator $\mathbf{A}(t, \mathbf{z}) : \mathbb{R}^3 \to \mathbb{R}^3$ such that $\mathbf{A}|_{T\Gamma_0} = J^{-1}D\Phi_t^n : T\Gamma_0 \to T\Gamma(t)$ and $\mathbf{A} : T\Gamma_0^{\perp} \to T\Gamma(t)^{\perp}$. We use the operator

$$\mathbf{A}(t,\mathbf{z})\mathbf{v} := J^{-1}(t,\mathbf{x})\mathbf{D}(t,\mathbf{z})\mathbf{v} + \mathbf{n}_{\Gamma(t)}(\mathbf{x})\mathbf{n}_{\Gamma_0}(\mathbf{z})\cdot\mathbf{v}, \quad \mathbf{v}\in\mathbb{R}^3, \ \mathbf{x}=\Phi_t^n(\mathbf{z}).$$
(3.4)

For $\mathbf{A}^{-1}(t, \mathbf{x}) := J(t, \mathbf{z})\mathbf{D}^{-1}(t, \mathbf{x}) + \mathbf{n}_{\Gamma_0}(\mathbf{z})\mathbf{n}_{\Gamma(t)}(\mathbf{x})^T$ it holds $\mathbf{A}^{-1}(t, \mathbf{x})\mathbf{A}(t, \mathbf{z}) = I_{\mathbb{R}^3}$. The matrices of \mathbf{A} and \mathbf{A}^{-1} in the standard basis are also denoted by \mathbf{A} and \mathbf{A}^{-1} , respectively. Note that det $\mathbf{A} = 1$ holds. We define the *surface Piola transform* $P_t : \mathbb{R}^3 \to \mathbb{R}^3$ by

$$(P_t \mathbf{v})(\mathbf{x}) := \mathbf{A}(t, \mathbf{z}) \, \mathbf{v}(\mathbf{z}), \quad \mathbf{z} \in \Gamma_0.$$
(3.5)

This operator maps tangential vectors on Γ_0 to tangential vectors on $\Gamma(t)$ and for tangential vectors \mathbf{v} it satisfies $\operatorname{div}_{\Gamma} P_t \mathbf{v} = 0$ a.e. on $\Gamma(t)$ iff $\operatorname{div}_{\Gamma} \mathbf{v} = 0$ a.e. on Γ_0 , cf. Ref. 41.

We need some regularity properties of **D**, **A** and **D**⁻¹, **A**⁻¹, which are collected in the following lemma. For a function $g \in C^1(\mathcal{S}_0)$ the maximum norm is $||g||_{C^1(\mathcal{S}_0)} := \max_{(t,\mathbf{z})\in\mathcal{S}_0}(|g(t,\mathbf{z})|+|\nabla_{S_0}g(t,\mathbf{z})|)$ and similarly for vector and matrix-valued functions as well as for such functions on \mathcal{S} .

Lemma 3.1. It holds that $\mathbf{D}, \mathbf{A} \in C^1(\mathcal{S}_0)^{3\times 3}, \mathbf{D}^{-1}, \mathbf{A}^{-1} \in C^1(\mathcal{S})^{3\times 3}$ and, in particular,

$$||J||_{C^{1}(\mathcal{S}_{0})} + ||\mathbf{D}||_{C^{1}(\mathcal{S}_{0})} + ||\mathbf{A}||_{C^{1}(\mathcal{S}_{0})} + ||J^{-1}||_{C^{1}(\mathcal{S})} + ||\mathbf{D}^{-1}||_{C^{1}(\mathcal{S})} + ||\mathbf{A}^{-1}||_{C^{1}(\mathcal{S})} \le C.$$
(3.6)

Proof. From (3.2), we know that $\Phi : (t, \mathbf{z}) \to (t, \Phi_t^n(\mathbf{z}))$ is in $C^2(\mathcal{S}_0, \mathcal{S})$ and hence $D\Phi \in C^1(T\mathcal{S}_0, T\mathcal{S})$. Moreover, $\mathbf{P}_{\mathcal{S}_0}$ is C^1 -smooth, so $(D\Phi)\mathbf{P}_{\mathcal{S}_0}$ is a C^1 smooth mapping with matrix representation

$$\begin{bmatrix} 1 & \mathbf{w}_N^T \\ 0 & \mathbf{D} \end{bmatrix} \in C^1(\mathcal{S}_0)^{4 \times 4}.$$
 (3.7)

Hence, $\mathbf{D} \in C^1(\mathcal{S}_0)^{3\times 3}$ and $J \in C^1(\mathcal{S}_0)$ hold. Combining this with $\mathbf{n}_{\Gamma_0}, \mathbf{n}_{\Gamma(\cdot)} \circ \Phi_{(\cdot)}^n \in C^1(\mathcal{S}_0)$, (3.4) and the property that \mathcal{S}_0 is closed, implies the bound in (3.6) for \mathbf{D} , J and \mathbf{A} . The mapping $\Phi : \mathcal{S}_0 \to \mathcal{S}$ is one-to-one. By the inverse mapping theorem the inverse Φ^{-1} is $C^2(\mathcal{S}, \mathcal{S}_0)$ and for its differential we have $D\Phi^{-1} \in C^1(T\mathcal{S}, T\mathcal{S}_0)$. The matrix of the C^1 smooth mapping $D\Phi^{-1}\mathbf{P}_{\mathcal{S}}$ is

$$\begin{bmatrix} 1 & -\mathbf{w}_N^T \mathbf{D}^{-1} \\ 0 & \mathbf{D}^{-1} \end{bmatrix} \in C^1(\mathcal{S})^{4 \times 4}.$$

This and $\mathbf{A}^{-1} = J\mathbf{D}^{-1} + \mathbf{n}_{\Gamma_0}\mathbf{n}_{\Gamma(\cdot)}^T(\Phi_{(\cdot)}^n)$ imply the desired bound for \mathbf{D}^{-1} , J^{-1} and \mathbf{A}^{-1} .

3.2. Evolving Hilbert spaces

For constructing suitable evolving Hilbert spaces, we first define *tangential* velocity spaces on $\Gamma(t)$. The notation $(\cdot, \cdot)_{0,t}$ and $\|\cdot\|_{0,t}$ is used for the canonical inner product and norm in $L^2(\Gamma(t))$. We need the Sobolev spaces of order one

$$H^{1}(t) := \{ \mathbf{v} \in H^{1}(\Gamma(t))^{3} \mid \mathbf{v} \cdot \mathbf{n} = 0 \text{ on}\Gamma(t) \},\$$

with the inner product $(\cdot, \cdot)_{1,t} := (\cdot, \cdot)_{0,t} + (D_{\Gamma} \cdot, D_{\Gamma} \cdot)_{0,t}$, and its closed subspace of divergence free tangential fields

$$V_1(t) := \{ \mathbf{v} \in H^1(t) \mid \operatorname{div}_{\Gamma} \mathbf{v} = 0 \text{ a.e. on } \Gamma(t) \}.$$

The space $V_0(t)$ is defined as closure of a space of smooth div-free tangential fields in the $L^2(\Gamma(t))$ norm:

$$V_0(t) := \overline{\mathcal{V}(t)}^{\|\cdot\|_{0,t}}, \quad \mathcal{V}(t) := \{ \mathbf{v} \in C^1(\Gamma(t))^3 \, | \, \mathbf{v} \cdot \mathbf{n} = 0, \, \operatorname{div}_{\Gamma} \mathbf{v} = 0 \text{ on } \Gamma(t) \}.$$

The space of smooth functions $\mathcal{V}(t)$ is dense not only in $V_0(t)$ but also in $V_1(t)$. Indeed, for any tangential velocity field $\mathbf{u} \in L^2(\Gamma(t))^3$ on the C^2 smooth surface $\Gamma(t)$ we have a Helmholtz decomposition $\mathbf{u} = \nabla_{\Gamma}\psi + \mathbf{n} \times (\nabla_{\Gamma}\phi) + \mathbf{h}$ with some $\psi, \phi \in H^1(\Gamma(t))$ and a harmonic field $\mathbf{h} \in C^1(\Gamma(t))^3$, see Ref. 35. For $\mathbf{u} \in V_1(t)$ we have $\psi = 0$ and the result follows from the density of C^2 -smooth functions in $H^1(\Gamma(t))$. Therefore, endowed with canonical scalar products, the spaces form a Gelfand triple $V_1(t) \hookrightarrow V_0(t) \hookrightarrow V_1(t)'$. We also have that the dense embedding $V_1(t) \hookrightarrow V_0(t)$ is compact. Here and later H' always denotes a dual of a Hilbert space H, and we adopt the common notation $H^{-1}(t)$ for $H^1(t)'$. For the space $L^2(t) := \{ \mathbf{v} \in L^2(\Gamma(t))^3 \mid \mathbf{v} \cdot \mathbf{n} = 0 \text{ a.e. on } \Gamma(t) \}$ we define a pushforward map $\phi_t : L^2(0) \to L^2(t)$, based on the Piola transform, by

$$(\phi_t \mathbf{v})(\mathbf{x}) = (P_t \mathbf{v})(\mathbf{x}) = \mathbf{A}(t, \mathbf{z})\mathbf{v}(\mathbf{z}), \quad \mathbf{v} \in L^2(0), \quad \mathbf{x} = \Phi_t^n(\mathbf{z}), \quad \mathbf{z} \in \Gamma_0.$$
(3.8)

The inverse map (pullback) is given by $(\phi_{-t}\mathbf{v})(\mathbf{z}) = \mathbf{A}^{-1}(t, \mathbf{x})\mathbf{v}(\mathbf{x}), \mathbf{v} \in L^2(t)$. Since $\mathbf{A} \in C^1(\mathcal{S}_0)^{3\times 3}$, the restriction of ϕ_t to $H^1(0)$ is a pushforward map from $H^1(0)$ to $H^1(t)$. Because ϕ_t is based on the Piola transform and thus conserves the solenoidal property, we also have that ϕ_t is a pushforward map from $V_0(0)$ to $V_0(t)$, and from $V_1(0)$ to $V_1(t)$. For this pushforward map we have for $\mathbf{v} \in H^1(0)$:

$$\begin{split} \|\phi_{t}\mathbf{v}\|_{1,t} &= \left(\|\mathbf{A}\mathbf{v}\circ\Phi_{-t}^{n}\|_{0,t}^{2} + \|D_{\Gamma}(\mathbf{A}\mathbf{v}\circ\Phi_{-t}^{n})\|_{0,t}^{2}\right)^{\frac{1}{2}} \\ &\leq \left(\|\mathbf{A}\|_{C(\Gamma(t))} + \|D_{\Gamma}\mathbf{A}\|_{C(\Gamma(t))}\right)\|\mathbf{v}\circ\Phi_{-t}^{n}\|_{0,t} + \|\mathbf{A}\|_{C(\Gamma(t))}\|D_{\Gamma}(\mathbf{v}\circ\Phi_{-t}^{n})\|_{0,t} \\ &\leq \left(\|\mathbf{A}\|_{C(\Gamma(t))} + \|D_{\Gamma}\mathbf{A}\|_{C(\Gamma(t))}\right)\|J\|_{C(\Gamma_{0})}^{\frac{1}{2}}\|\mathbf{v}\|_{0,0} \\ &+ \|\mathbf{A}\|_{C(\Gamma(t))}\|\mathbf{D}^{-1}\|_{C(\Gamma(t))}\|J\|_{C(\Gamma_{0})}^{\frac{1}{2}}\|D_{\Gamma}\mathbf{v}\|_{0,0}. \end{split}$$

The result (3.6) implies that the norms $\|\mathbf{A}\|_{C(\Gamma(t))}$, $\|\mathbf{D}^{-1}\|_{C(\Gamma(t))}$, $\|D_{\Gamma}\mathbf{A}\|_{C(\Gamma(t))}$, $\|J\|_{C(\Gamma_0)}$ are uniformly bounded in t and thus

$$\sup_{t \in [0,T]} \|\phi_t \mathbf{v}\|_{1,t} \le C \|\mathbf{v}\|_{1,0}$$

holds with some C independent of $\mathbf{v} \in H^1(0)$. With similar arguments one easily shows that $\|\phi_{-t}\mathbf{v}\|_{1,0} \leq C \|\mathbf{v}\|_{1,t}$ holds for all $\mathbf{v} \in H^1(t)$, with C independent of \mathbf{v} and t. These bounds remain obviously true if $H^1(0)$, $H^1(t)$ and the corresponding norms are replaced by $V_0(0)$, $V_0(t)$ and the corresponding norms. Using (3.6) one shows that the maps $t \to \|\phi_t \mathbf{v}\|_{1,t}$ and $t \to \|\phi_t \mathbf{v}\|_{0,t}$ are continuous. These properties imply that $\{V_0(t), \phi_t\}_{t \in [0,T]}, \{H^1(t), \phi_t\}_{t \in [0,T]}, \text{ and } \{V_1(t), \phi_t\}_{t \in [0,T]}$ are "compatible pairs" in the sense of Definition 2.4 in Ref. 1. This compatibility structure induces some useful properties of the evolving spaces defined as follows:

$$L^{2}_{V_{1}} := \left\{ \mathbf{v} : [0,T] \to \bigcup_{t \in [0,T]} \{t\} \times V_{1}(t), t \to (t, \bar{\mathbf{v}}(t)) \, | \, \phi_{-(\cdot)} \bar{\mathbf{v}}(\cdot) \in L^{2}(0,T; V_{1}(0)) \right\},$$

$$L^{2}_{V'_{1}} := \left\{ \mathbf{g} : [0,T] \to \bigcup_{t \in [0,T]} \{t\} \times V_{1}(t)', t \to (t, \bar{\mathbf{g}}(t)) \mid \phi^{*}_{(\cdot)} \bar{\mathbf{g}}(\cdot) \in L^{2}(0,T; V_{1}(0)') \right\},\$$

$$L^2_{V_0} := \left\{ \mathbf{v} : [0,T] \to \bigcup_{t \in [0,T]} \{t\} \times V_0(t), t \to (t, \bar{\mathbf{v}}(t)) \, | \, \phi_{-(\cdot)} \bar{\mathbf{v}}(\cdot) \in L^2(0,T;V_0(0)) \right\},$$

where ϕ_t^* is the dual of ϕ_t . We shall also need the spaces $L_{V_0}^{\infty}$, $L_{H^1}^2$, $L_{H^{-1}}^2$, which are defined analogously, and the spaces of smooth space-time functions

$$\mathcal{D} = \{ \mathbf{v} \in L^2_{V_1} \, | \, \phi_{-(\cdot)} \bar{\mathbf{v}}(\cdot) \in C^{\infty}([0,T]; \mathcal{V}(0)) \}, \\ \mathcal{D}_0 = \{ \mathbf{v} \in L^2_{V_1} \, | \, \phi_{-(\cdot)} \mathbf{v}(\cdot) \in C^{\infty}_0((0,T); \mathcal{V}(0)) \}.$$
(3.9)

Note that functions in \mathcal{D}_0 have zero traces on ∂S . With a slight abuse of notation we identify $\bar{\mathbf{v}}(t)$ with $\mathbf{v}(t) = (t, \bar{\mathbf{v}}(t))$.

In Ref. 1, it is shown that if $V_1(0) \hookrightarrow V_0(0) \hookrightarrow V_1(0)'$ is a Gelfand triple, with a compact embedding $V_1(0) \hookrightarrow V_0(0)$, and both $\{V_0(t), \phi_t\}_{t \in [0,T]}$ and $\{V_1(t), \phi_t\}_{t \in [0,T]}$ are compatible pairs, then the *L*-spaces inherit certain properties of the standard Bochner spaces. In particular, cf. Sec. 2 in Ref. 1, the spaces $L^2_{V_1}$ and $L^2_{V_0}$ with

$$(\mathbf{u}, \mathbf{v})_{\mathbf{1}} = \int_0^T (\mathbf{u}(t), \mathbf{v}(t))_{1,t} dt, \quad (\mathbf{u}, \mathbf{v})_{\mathbf{0}} = \int_0^T (\mathbf{u}(t), \mathbf{v}(t))_{0,t} dt,$$

are separable Hilbert spaces, homeomorphic to $L^2(0,T;V_1(0))$ and $L^2(0,T;V_0(0))$, respectively. Furthermore, the embedding $L^2_{V_1} \hookrightarrow L^2_{V_0}$ is dense and compact, the dual space $(L^2_{V_1})'$ is isometrically isomorphic to $L^2_{V_1}$, and

$$L^2_{V_1} \hookrightarrow L^2_{V_0} \hookrightarrow L^2_{V_1'}$$

is a Gelfand triple. The space \mathcal{D}_0 is dense in $L^2_{V_1}$ and so \mathcal{D}_0 is dense also in $L^2_{V_0}$. By the same arguments $L^2_{H^1}$ is also a Hilbert space with inner product $(\cdot, \cdot)_1$. The subspace of smooth functions

$$\widehat{\mathcal{D}}_0 = \{ \mathbf{v} \in L^2_{H^1} \mid \phi_{-(\cdot)} \mathbf{v}(\cdot) \in C^\infty_0((0,T); C^1(\Gamma_0, T\Gamma_0)) \}$$
(3.10)

is dense in $L^2_{H^1}$ and $(L^2_{H^1})' \simeq L^2_{H^{-1}}$ holds. Note that $L^2_{V_1}$ is a closed subspace of $L^2_{H^1}$ and that functions in \mathcal{D}_0 are solenoidal and functions in $\widehat{\mathcal{D}}_0$ are not necessarily solenoidal.

3.3. Some uniform inequalities

We need to establish several basic inequalities on $\Gamma(t)$ with constants uniformly bounded in t.

We first consider a Korn inequality. Recall that the estimate

$$\|\mathbf{v}\|_{1,t} \le c(\|\mathbf{v}\|_{0,t} + \|E_s(\mathbf{v})\|_{0,t}) \quad \text{for all } \mathbf{v} \in H^1(t)$$
(3.11)

holds with a constant c = c(t) that depends on smoothness properties of $\Gamma(t)$, cf. Ref. 16. In the next lemma, we show that the constant can be taken such that $\max_{t \in [0,T]} c(t) < \infty$ holds.

Lemma 3.2. The constant c in (3.11) can be chosen finite and independent of t.

Proof. Fix any $t \in [0,T]$ and $\mathbf{v} \in H^1(t)$. Define $\mathbf{u} = \mathbf{D}^T \mathbf{v} \circ \Phi_t^n \in H^1(0)$. Below, for the 3-tensor $\mathbb{T} = D_{\Gamma} \mathbf{D}^T$ and $\mathbf{v} \in \mathbb{R}^3$, $\mathbb{T} \mathbf{v} \in \mathbb{R}^{3\times 3}$ is the second-mode tensor-vector product. With the help of (3.3) one computes

$$\begin{aligned} \nabla_{\Gamma_0} \mathbf{u}(\mathbf{z}) &= \mathbf{P}(\mathbf{z})[D_{\Gamma_0} \mathbf{u}(\mathbf{z})] = \mathbf{P}(\mathbf{z})[D_{\Gamma_0} (\mathbf{D}^T \mathbf{v} \circ \Phi_t^n(\mathbf{z}))] \\ &= \mathbf{P}(\mathbf{z}) \mathbb{T} \mathbf{v}(\mathbf{x}) + \mathbf{P}(\mathbf{z}) \mathbf{D}^T D_{\Gamma} \mathbf{v}(\mathbf{x}) \mathbf{D}, \\ &= \mathbf{P}(\mathbf{z}) \mathbb{T} \mathbf{v}(\mathbf{x}) + \mathbf{D}^T \nabla_{\Gamma} \mathbf{v}(\mathbf{x}) \mathbf{D}, \quad \mathbf{x} = \Phi_t^n(\mathbf{z}). \end{aligned}$$

The latter inequality holds since $\mathbf{D}^T \mathbf{P}(\mathbf{x}) = \mathbf{P}(\mathbf{z}) \mathbf{D}^T$. From this we find

$$E_s(\mathbf{u})(\mathbf{z}) = \frac{1}{2} (\mathbf{P}(\mathbf{z}) \mathbb{T} \mathbf{v}(\mathbf{x}) + (\mathbb{T} \mathbf{v}(\mathbf{x}))^T \mathbf{P}(\mathbf{z})) + \mathbf{D}^T E_s(\mathbf{v})(\mathbf{x}) \mathbf{D}.$$
 (3.12)

With the help of $D_{\Gamma} \mathbf{v} = D_{\Gamma}(\mathbf{D}^{-T} \mathbf{u} \circ \Phi_{-t}^n)$, (3.6), (3.11) applied for t = 0, and (3.12) we estimate

$$\begin{split} \|D_{\Gamma}\mathbf{v}\|_{0,t} &= \|J^{\frac{1}{2}}([D_{\Gamma}\mathbf{D}^{-T}]\mathbf{u} + \mathbf{D}^{-T}(D_{\Gamma}\mathbf{u})\mathbf{D}^{-1})\|_{0,0} \\ &\leq C(\|\mathbf{u}\|_{0,0} + \|D_{\Gamma}\mathbf{u}\|_{0,0}) \leq C(\|\mathbf{u}\|_{0,0} + \|E_{s}(\mathbf{u})\|_{0,0}) \\ &= C\left(\|J^{-\frac{1}{2}}\mathbf{D}^{T}\mathbf{v}\|_{0,t} + \left\|J^{-\frac{1}{2}}\left(\frac{1}{2}(\mathbf{P}\mathbb{T}\mathbf{v} + (\mathbb{T}\mathbf{v})^{T}\mathbf{P}) + \mathbf{D}^{T}E_{s}(\mathbf{v})\mathbf{D}\right)\right\|_{0,t}\right) \\ &\leq C(\|J^{-\frac{1}{2}}\mathbf{D}^{T}\mathbf{v}\|_{0,t} + \|J^{-\frac{1}{2}}\mathbf{P}\mathbb{T}\mathbf{v}\|_{0,t} + \|J^{-\frac{1}{2}}\mathbf{D}^{T}E_{s}(\mathbf{v})\mathbf{D}\|_{0,t}) \\ &\leq C(\|\mathbf{v}\|_{0,t} + \|E_{s}(\mathbf{v})\|_{0,t}) \end{split}$$

with some C independent of $t \in [0, T]$ and **v**.

The following inf-sup estimate holds¹⁶: there exists c(t) > 0 such that

$$\|\nabla_{\Gamma}\pi\|_{H^{-1}(t)} := \sup_{0 \neq \mathbf{v} \in H^{1}(t)} \frac{\int_{\Gamma(t)} \pi \operatorname{div}_{\Gamma} \mathbf{v} \, ds}{\|\mathbf{v}\|_{1,t}} \ge c(t) \|\pi\|_{0,t}, \quad \forall \, \pi \in L^{2}(\Gamma(t)),$$
(3.13)

with $\int_{\Gamma(t)} \pi = 0$. A uniformity result for this inf-sup constant is derived in the following lemma.

Lemma 3.3. The constant c(t) in (3.13) can be taken such that $\inf_{t \in [0,T]} c(t) > 0$ holds.

Proof. We use a similar approach as in the proof of the previous lemma, and derive an estimate on $\Gamma(t)$ by pulling forward the result on Γ_0 . We use the pullforward ϕ_t that is based on the Piola transform and satisfies $\operatorname{div}_{\Gamma(t)}(\phi_t \mathbf{w})(\mathbf{x}) = J^{-1} \operatorname{div}_{\Gamma_0} \mathbf{w}(\mathbf{z})$, $\mathbf{z} \in \Gamma_0, \mathbf{x} = \Phi_t^n(\mathbf{z}) \in \Gamma(t)$. Take $\mathbf{v} \in H^1(t)$ and $\pi \in L^2(\Gamma(t))$ with $\int_{\Gamma(t)} \pi = 0$. Define $c := -|\Gamma_0|^{-1} \int_{\Gamma_0} \pi \circ \Phi_t^n ds$ and $\mathbf{w} := \phi_{-t} \mathbf{v} \in H^1(0)$. Note that $\|\mathbf{v}\|_{1,t} \leq C \|\mathbf{w}\|_{1,0}$ with a constant C uniformly bounded in $t \in [0,T]$ (compatibility property). We have

$$(\pi, \operatorname{div}_{\Gamma(t)}\mathbf{v})_{0,t} = (\pi, \operatorname{div}_{\Gamma(t)}(\phi_t \mathbf{w}))_{0,t} = (\pi \circ \Phi_t^n, \operatorname{div}_{\Gamma_0} \mathbf{w})_{0,0}.$$

Using this and the result (3.13) for t = 0 we get

$$C \sup_{0 \neq \mathbf{v} \in H^{1}(t)} \frac{\int_{\Gamma(t)} \pi \operatorname{div}_{\Gamma} \mathbf{v} \, ds}{\|\mathbf{v}\|_{1,t}} \ge \sup_{0 \neq \mathbf{w} \in H^{1}(0)} \frac{\int_{\Gamma_{0}} (\pi \circ \Phi_{t}^{n}) \operatorname{div}_{\Gamma_{0}} \mathbf{w} \, ds}{\|\mathbf{w}\|_{1,0}}$$
$$\ge c(0) \|\pi \circ \Phi_{t}^{n} + c\|_{0,0} \ge c(0) \|J^{\frac{1}{2}}\|_{C(\mathcal{S})}^{-1} \|\pi + c\|_{0,t}$$

$$= c(0) \|J^{\frac{1}{2}}\|_{C(\mathcal{S})}^{-1}(\|\pi\|_{0,t} + c|\Gamma(t)|^{\frac{1}{2}})$$

$$\geq c(0) \|J^{\frac{1}{2}}\|_{C(\mathcal{S})}^{-1}\|\pi\|_{0,t},$$

which yields a *t*-independent strictly positive lower bound for c(t) in (3.13).

We now derive a uniform interpolation estimate.

Lemma 3.4. The interpolation inequality (Ladyzhenskaya's inequality)

$$\|\mathbf{v}\|_{L^{4}(\Gamma(t))} \leq C \|\mathbf{v}\|_{0,t}^{\frac{1}{2}} \|\mathbf{v}\|_{1,t}^{\frac{1}{2}}, \quad \mathbf{v} \in H^{1}(t).$$
(3.14)

holds with a constant $C < \infty$ independent of t.

Proof. Consider $v \in H^1(\Gamma(t))$ and let $\hat{v} = v \circ \Phi_{-t}^n$. For a compact Riemann manifold Γ_0 , the estimate (II.38) from Ref. 3 yields $\|\hat{v}\|_{L^4(\Gamma_0)} \leq C \|\hat{v}\|_{0,0}^{\frac{1}{2}} \|\hat{v}\|_{1,0}^{\frac{1}{2}}$. An examination of the proof shows that the estimate remains true if Γ_0 is a C^2 compact manifold. With the help of this estimate applied component-wise and (3.6) we calculate for $\mathbf{v} \in H^1(t)$

$$\begin{aligned} \|\mathbf{v}\|_{L^{4}(\Gamma(t))} &= \|\hat{\mathbf{v}}J^{\frac{1}{4}}\|_{L^{4}(\Gamma_{0})} \leq C \|J\|_{C(\mathcal{S}_{0})}^{\frac{1}{4}}\|\hat{\mathbf{v}}\|_{L^{4}(\Gamma_{0})} \leq C \|\hat{\mathbf{v}}\|_{0,0}^{\frac{1}{2}} \|\hat{\mathbf{v}}\|_{1,0}^{\frac{1}{2}} \\ &\leq C \|J^{-\frac{1}{2}}\mathbf{v}\|_{0,t}^{\frac{1}{2}} \left(\|J^{-\frac{1}{2}}\mathbf{v}\|_{0,t} + \|J^{-\frac{1}{2}}\mathbf{D}^{-T}\nabla_{\Gamma}\mathbf{v}\|_{0,t}\right)^{\frac{1}{2}} \\ &\leq C \|\mathbf{v}\|_{0,t}^{\frac{1}{2}} \left(\|\mathbf{v}\|_{0,t} + \|\nabla_{\Gamma}\mathbf{v}\|_{0,t}\right)^{\frac{1}{2}} \leq C \|\mathbf{v}\|_{0,t}^{\frac{1}{2}} \|\mathbf{v}\|_{1,t}^{\frac{1}{2}}, \end{aligned}$$

with some C independent of t.

For $\boldsymbol{\xi} \in H^1(t)$ consider the Helmholtz decomposition (see e.g. Ref. 35)

$$\boldsymbol{\xi} = \boldsymbol{\xi}_1 + \boldsymbol{\xi}_2, \quad \text{with } \boldsymbol{\xi}_1 = \nabla_{\Gamma} \phi, \quad \phi \in H^1(\Gamma(t)) \quad \text{and} \quad \boldsymbol{\xi}_2 \in V_1(t).$$
(3.15)

Lemma 3.5. For $\boldsymbol{\xi}_i$ as in (3.15) we have $\boldsymbol{\xi}_i \in H^1(t)$ and $\|\boldsymbol{\xi}_i\|_{1,t} \leq C\|\boldsymbol{\xi}\|_{1,t}$, i = 1, 2, with a constant C finite and independent of t.

Proof. Due to the L^2 orthogonality of the Helmholtz decomposition we have $\|\boldsymbol{\xi}_1\|_{0,t}^2 + \|\boldsymbol{\xi}_2\|_{0,t}^2 = \|\boldsymbol{\xi}\|_{0,t}^2$. Also note that $\operatorname{div}_{\Gamma}\boldsymbol{\xi}_2 = 0$, $\operatorname{div}_{\Gamma}\boldsymbol{\xi} = \operatorname{div}_{\Gamma}\boldsymbol{\xi}_1$, $\operatorname{curl}_{\Gamma}\boldsymbol{\xi}_1 = 0$, $\operatorname{curl}_{\Gamma}\boldsymbol{\xi} = \operatorname{curl}_{\Gamma}\boldsymbol{\xi}_2$. Furthermore on $H^1(t)$ we have the norm equivalence $\|\mathbf{u}\|_{1,t} \sim \|\mathbf{u}\|_{0,t} + \|\operatorname{div}_{\Gamma}\mathbf{u}\|_{0,t} + \|\operatorname{curl}_{\Gamma}\mathbf{u}\|_{0,t}$. A *t*-dependence in the constants in this norm equivalence enters only through the Gaussian curvature of $\Gamma(t)$, cf. Ref. 35. Due to the smoothness property $\mathcal{S} \in C^3$ the Gaussian curvature is uniformly bounded on \mathcal{S} and thus the constants in this norm equivalence can be taken independent of *t*. Using these results we get

$$\begin{split} \|\boldsymbol{\xi}_1\|_{1,t} &\leq C(\|\boldsymbol{\xi}_1\|_{0,t} + \|\operatorname{div}_{\Gamma}\boldsymbol{\xi}_1\|_{0,t} + \|\operatorname{curl}_{\Gamma}\boldsymbol{\xi}_1\|_{0,t}) \\ &= C(\|\boldsymbol{\xi}_1\|_{0,t} + \|\operatorname{div}_{\Gamma}\boldsymbol{\xi}\|_{0,t}) \leq C\|\boldsymbol{\xi}\|_{1,t}, \end{split}$$

and by similar arguments $\|\boldsymbol{\xi}_2\|_{1,t} \leq C \|\boldsymbol{\xi}\|_{1,t}$ with a constant *C* uniformly bounded in *t*.

3.4. Solution space

In this section, we introduce a subspace of $L_{V_1}^2$ consisting of functions for which a suitable weak normal time derivative exists. This space will be the solution space in the weak formulation of TSNSE.

We recall the Leibniz rule

$$\frac{d}{dt} \int_{\Gamma(t)} v \, ds = \int_{\Gamma(t)} (\partial^{\circ} v + v \operatorname{div}_{\Gamma} \mathbf{w}_N) \, ds = \int_{\Gamma(t)} (\partial^{\circ} v + v w_N \kappa) \, ds,$$

Thus for velocity fields $\mathbf{v}, \mathbf{u} \in C^1(\mathcal{S})$ we get

$$\frac{d}{dt} \int_{\Gamma(t)} \mathbf{v} \cdot \mathbf{u} \, ds = \int_{\Gamma(t)} (\partial^{\circ} (\mathbf{v} \cdot \mathbf{u}) + (\mathbf{v} \cdot \mathbf{u}) w_N \kappa) \, ds.$$
(3.16)

This implies the integration by parts identity

$$\int_{0}^{T} \int_{\Gamma(t)} (\partial^{\circ} \mathbf{v} \cdot \mathbf{u} + \mathbf{v} \cdot \partial^{\circ} \mathbf{u} + (\mathbf{v} \cdot \mathbf{u}) w_{N} \kappa) \, ds \, dt$$
$$= \int_{\Gamma(T)} \mathbf{v} \cdot \mathbf{u} \, ds - \int_{\Gamma_{0}} \mathbf{v} \cdot \mathbf{u} \, ds, \quad \mathbf{v}, \mathbf{u} \in C^{1}(\mathcal{S})^{3}.$$
(3.17)

Based on this we define for $\mathbf{v} \in L^2_{H^1}$ the normal time derivative as the functional $\partial^{\circ} \mathbf{v}$:

$$\langle \partial^{\circ} \mathbf{v}, \boldsymbol{\xi} \rangle := -\int_{0}^{T} \int_{\Gamma(t)} (\mathbf{v} \cdot \partial^{\circ} \boldsymbol{\xi} + (\mathbf{v} \cdot \boldsymbol{\xi}) w_{N} \kappa) \, ds \, dt, \quad \boldsymbol{\xi} \in \widehat{\mathcal{D}}_{0}.$$
(3.18)

Note that functions in $\widehat{\mathcal{D}}_0$ are not necessarily solenoidal, cf. (3.10). Restricting now to $L^2_{V_1} \subset L^2_{H^1}$, assume $\mathbf{v} \in L^2_{V_1}$ is such that

$$\|\partial^{\circ}\mathbf{v}\|_{(L^{2}_{V_{1}})'} := \sup_{\boldsymbol{\xi}\in\mathcal{D}_{0}}rac{\langle\partial^{\circ}\mathbf{v},\boldsymbol{\xi}
angle}{\|\boldsymbol{\xi}\|_{1}}$$

is bounded. Since \mathcal{D}_0 is dense in $L^2_{V_1}$, $\partial^{\circ} \mathbf{v}$ can then be extended to a bounded linear functional on $L^2_{V_1}$. We use $(L^2_{V_1})' \cong L^2_{V'_1}$ and introduce the space

$$\mathbf{W}(V_1, V_1') = \{ \mathbf{v} \in L^2_{V_1} \mid \partial^{\circ} \mathbf{v} \in L^2_{V_1'} \}, \text{ with}$$
$$(\mathbf{v}, \mathbf{u})_W := \int_0^T (\mathbf{v}(t), \mathbf{u}(t))_{1,t} + (\partial^{\circ} \mathbf{v}(t), \partial^{\circ} \mathbf{u}(t))_{V_1(t)'} dt.$$

This space is used as solution space in the weak formulation of TSNSE below. In the remainder of this section, we derive certain useful properties of this space. For this it will be helpful to introduce in addition to the Lagrangian derivatives $\dot{\mathbf{v}}$ (material derivative) and $\partial^{\circ}\mathbf{v}$ (normal time derivative) one other Lagrangian derivative, which is based on the pushforward operator ϕ_t :

$$\partial^* \mathbf{v}(t) := \phi_t \left(\frac{d}{dt} \phi_{-t} \mathbf{v}(t) \right), \quad \mathbf{v} \in \mathcal{D}.$$
(3.19)

The reason that we introduce the ∂^* derivative is, that it is the same as the one used in the general framework in Ref. 1 and we can use results derived in that paper. Note that the ∂^* derivative is defined for tangential flow fields and based on the Piola transform implying

$$\mathbf{n} \cdot \partial^* \mathbf{v} = 0$$
 and $\operatorname{div}_{\Gamma} \partial^* \mathbf{v} = 0$ for $\mathbf{v} \in \mathcal{D}$. (3.20)

We now derive relations between the derivatives ∂^* and ∂° .

Lemma 3.6. For $\mathbf{v} \in \mathcal{D}$ the following holds:

$$\partial^{\circ} \mathbf{v} = \partial^* \mathbf{v} - \mathbf{A} (\partial^{\circ} \mathbf{A}^{-1}) \mathbf{v}, \qquad (3.21)$$

$$\mathbf{P}\partial^{\circ}\mathbf{v} = \partial^{*}\mathbf{v} - \mathbf{A}\mathbf{P}(\partial^{\circ}\mathbf{A}^{-1})\mathbf{v}.$$
(3.22)

Proof. Using the definitions of the pushforward and pullback mappings we compute

$$\begin{split} \phi_t \left(\frac{d}{dt} \phi_{-t} \mathbf{v}(t) \right) (\mathbf{x}) &= \mathbf{A}(t, \mathbf{z}) \frac{d}{dt} \left[\mathbf{A}^{-1}(t, \Phi_t^n(\mathbf{z})) \mathbf{v}(t, \Phi_t^n(\mathbf{z})) \right] \\ &= \mathbf{A}(t, \mathbf{z}) (\partial^{\circ} \mathbf{A}^{-1}(t, \mathbf{x}) \mathbf{v}(t, \mathbf{x}) + \mathbf{A}^{-1}(t, \mathbf{x}) \partial^{\circ} \mathbf{v}(t, \mathbf{x})) \\ &= \mathbf{A}(t, \mathbf{z}) \partial^{\circ} \mathbf{A}^{-1}(t, \mathbf{x}) \mathbf{v}(t, \mathbf{x}) + \partial^{\circ} \mathbf{v}(t, \mathbf{x}), \end{split}$$

which yields the result (3.21). The result (3.22) follows from (3.21) using $\mathbf{P}\partial^*\mathbf{v} = \partial^*\mathbf{v}$ and $\mathbf{P}\mathbf{A} = \mathbf{P}(\mathbf{x})\mathbf{A}(t, \mathbf{z}) = \mathbf{A}(t, \mathbf{z})\mathbf{P}(\mathbf{z})$.

From (3.21) we obtain the identity

$$(\partial^* \mathbf{v}, \boldsymbol{\xi})_{\mathbf{0}} = (\partial^\circ \mathbf{v}, \boldsymbol{\xi})_{\mathbf{0}} + (\mathbf{C}\mathbf{v}, \boldsymbol{\xi})_{\mathbf{0}}, \quad \forall \, \mathbf{v} \in \mathcal{D}, \ \boldsymbol{\xi} \in \widehat{\mathcal{D}}_0,$$

with $\mathbf{C} := \mathbf{AP}(\partial^{\circ} \mathbf{A}^{-1})$. Based on this, we define $\partial^* \mathbf{v}$ for $\mathbf{v} \in L^2_{V_1}$ as the functional

$$\langle \partial^* \mathbf{v}, \boldsymbol{\xi} \rangle := \langle \partial^\circ \mathbf{v}, \boldsymbol{\xi} \rangle + (\mathbf{C} \mathbf{v}, \boldsymbol{\xi})_{\mathbf{0}}, \quad \boldsymbol{\xi} \in \widehat{\mathcal{D}}_0.$$
 (3.23)

with $\langle \partial^{\circ} \mathbf{v}, \boldsymbol{\xi} \rangle$ defined in (3.18). The density of $\widehat{\mathcal{D}}_0$ in $L^2_{H^1}$ and of $\mathcal{D}_0 \subset \widehat{\mathcal{D}}_0$ in $L^2_{V_1}$ allows us to define $\partial^* \mathbf{v}$ as an element of $L^2_{H^{-1}}$ and $L^2_{V'_1}$, respectively. The following result holds:

$$\partial^* \mathbf{v} \in L^2_{V'_1} \Leftrightarrow \partial^* \mathbf{v} \in L^2_{H^{-1}}, \quad \mathbf{v} \in L^2_{V_1}.$$
(3.24)

Implication " \Leftarrow " in (3.24) is trivial since $V_1 \subset H^1$. To see " \Rightarrow ", consider any $\mathbf{v} \in \mathcal{D}$ and $\boldsymbol{\xi} \in L^2_{H^1}$ together with its Helmholtz decomposition $\boldsymbol{\xi} = \nabla_{\Gamma}\phi + \boldsymbol{\xi}_2$, cf. (3.15). Thanks to Lemma 3.5 we get $\nabla_{\Gamma}\phi \in L^2_{H^1}$, $\boldsymbol{\xi}_2 \in L^2_{V_1}$ and $\|\nabla_{\Gamma}\phi\|_1 + \|\boldsymbol{\xi}_2\|_1 \leq C\|\boldsymbol{\xi}\|_1$. Since $\boldsymbol{\xi}_2 \in L^2_{V_1}$ we have

$$|\langle \partial^* \mathbf{v}, \boldsymbol{\xi}_2 \rangle| \le \|\partial^* \mathbf{v}\|_{L^2_{V_1'}} \|\boldsymbol{\xi}_2\|_{\mathbf{1}} \le C \|\partial^* \mathbf{v}\|_{L^2_{V_1'}} \|\boldsymbol{\xi}\|_{\mathbf{1}}, \qquad (3.25)$$

while for the other component we get employing (3.20)

$$\langle \partial^* \mathbf{v}, \nabla_{\Gamma} \phi \rangle = (\partial^* \mathbf{v}, \nabla_{\Gamma} \phi)_{\mathbf{0}} = -(\operatorname{div}_{\Gamma} \partial^* \mathbf{v}, \phi)_{\mathbf{0}} = 0.$$
(3.26)

We thus conclude $|\langle \partial^* \mathbf{v}, \boldsymbol{\xi} \rangle| \leq C \|\partial^* \mathbf{v}\|_{L^2_{V'_1}} \|\boldsymbol{\xi}\|_1$ for all $\mathbf{v} \in \mathcal{D}$ and $\boldsymbol{\xi} \in L^2_{H^1}$. The result in (3.24) follows from the density of \mathcal{D} in $L^2_{V_1}$.

We are now ready to prove the following result.

Lemma 3.7. The space $\mathbf{W}(V_1, V'_1)$ is a Hilbert space and \mathcal{D} is dense in $\mathbf{W}(V_1, V'_1)$. For any $\mathbf{v} \in \mathbf{W}(V_1, V'_1)$ and $t \in [0, T]$, $\mathbf{v}(t)$ is well defined as an element of $V_0(t)$ and it holds

$$\sup_{t \in [0,T]} \|\mathbf{v}(t)\|_{0,t} \le C \|\mathbf{v}\|_W.$$

Proof. The idea of the proof is to relate the space $\mathbf{W}(V_1, V_1')$ to the space $\mathbf{W}_*(V_1, V_1') := \{ \mathbf{v} \in L^2_{V_1} | \partial^* \mathbf{v} \in L^2_{V_1'} \}$, with $\| \cdot \|_{W_*} = (\| \cdot \|_1^2 + \| \partial^* \cdot \|_{L^2_{V_1'}}^2)^{\frac{1}{2}}$, and to show that the latter is homeomorphic to a standard Bochner space for \mathcal{S}_0 . Lemma 3.1 ensures $\mathbf{C} \in C(\mathcal{S})^{3 \times 3}$ and thus from (3.23) we obtain

$$|\langle \partial^* \mathbf{v}, \boldsymbol{\xi} \rangle| - c \|\mathbf{v}\|_{\mathbf{0}} \|\boldsymbol{\xi}\|_{\mathbf{0}} \le |\langle \partial^\circ \mathbf{v}, \boldsymbol{\xi} \rangle| \le |\langle \partial^* \mathbf{v}, \boldsymbol{\xi} \rangle| + c \|\mathbf{v}\|_{\mathbf{0}} \|\boldsymbol{\xi}\|_{\mathbf{0}}.$$

Therefore, $\partial^* \mathbf{v}$ is a linear bounded functional on $L^2_{V_1}$ iff $\partial^{\circ} \mathbf{v}$ has this property. We conclude $\mathbf{v} \in \mathbf{W}(V_1, V'_1)$ iff $\mathbf{v} \in \mathbf{W}_*(V_1, V'_1)$. Moreover, the above inequalities, definition of the $\mathbf{W}(V_1, V'_1)$ -norm, $\mathbf{W}_*(V_1, V'_1)$ -norm and $L^2_{V_1} \hookrightarrow L^2_{V_0}$ yield

$$c \|\mathbf{v}\|_W \le \|\mathbf{v}\|_{W_*} \le C \|\mathbf{v}\|_W,$$

with constants 0 < c and $C < +\infty$ independent of $\mathbf{v} \in \mathbf{W}(V_1, V_1')$ and so $\mathbf{W}(V_1, V_1') = \mathbf{W}_*(V_1, V_1')$ algebraically and topologically. Thus, it is sufficient to check the claims of the lemma for $\mathbf{W}_*(V_1, V_1')$. For the latter we apply results from Ref. 1, more specifically, Corollary 2.32 ($\mathbf{W}_*(V_1, V_1')$ is a Hilbert space), Lemma 2.35 (continuous embedding $\mathbf{W}_*(V_1, V_1') \hookrightarrow C([0, T]; V_0(0))$) and Lemma 2.38 (density of smooth functions). For these results to hold one has to verify Assumption 2.31 in Ref. 1, which requires the mapping $\mathbf{v} \to \phi_{-(.)}\mathbf{v}$ to be a homeomorphism between $\mathbf{W}_*(V_1, V_1')$ and $\mathcal{W}(V_1(0), V_1(0)')$, the standard Bochner space

$$\mathcal{W}(V_1(0), V_1(0)') = \{ \mathbf{v} \in L^2((0, T), V_1(0)) \mid \partial_t \mathbf{v} \in L^2((0, T), V_1(0)') \}.$$

It remains to check this homeomorphism property. We already derived the norm equivalence $\|\mathbf{v}\|_{\mathbf{1}} \simeq \|\phi_{-(.)}\mathbf{v}\|_{L^2([0,T],V_1(0))}$, cf. Sec. 3.2. To relate the norms $\|\partial_t \phi_{-(.)}\mathbf{v}\|_{L^2((0,T),V_1(0)')}$ and $\|\partial^*\mathbf{v}\|_{L^2_{V_1'}}$ we consider the following equalities for $\mathbf{v} \in L^2_{V_1}$, $\boldsymbol{\xi} \in \mathcal{D}_0$, $\tilde{\boldsymbol{\xi}} = \phi_{-(.)}\boldsymbol{\xi} \in C^\infty_0((0,T),\mathcal{V}(0))$ and $\mathbf{A}^{-1}(t,\cdot) : T\Gamma(t) \to T\Gamma_0$:

$$\begin{aligned} \langle \partial_t(\phi_{-(\cdot)}\mathbf{v}), \hat{\boldsymbol{\xi}} \rangle &= (\phi_{-(\cdot)}\mathbf{v}, \partial_t \hat{\boldsymbol{\xi}})_{L^2(\mathcal{S}_0)} = (\mathbf{A}^{-1}\mathbf{v}(\cdot, \Phi_t^n(\cdot)), \partial_t(\phi_{-(\cdot)}\boldsymbol{\xi}))_{L^2(\mathcal{S}_0)} \\ &= (\mathbf{A}^{-T}\mathbf{A}^{-1}\mathbf{v}(\cdot, \Phi_t^n(\cdot)), \mathbf{A}\partial_t(\phi_{-(\cdot)}\boldsymbol{\xi}))_{L^2(\mathcal{S}_0)} \\ &= (J^{-1}(\mathbf{A}\mathbf{A}^T)^{-1}\mathbf{v}, \partial^*\boldsymbol{\xi})_{\mathbf{0}} \\ &= (\mathbf{v}, \partial^*(\mathbf{T}\boldsymbol{\xi}) - (\partial^*\mathbf{T})\boldsymbol{\xi})_{\mathbf{0}} \quad \text{with } \mathbf{T} := J^{-1}(\mathbf{A}\mathbf{A}^T)^{-1}. \end{aligned}$$

Note that $\phi := \mathbf{T}\boldsymbol{\xi}$ is smooth and has zero trace on ∂S . Using this, (3.22), (3.18) and (3.23), we have

$$\begin{aligned} (\mathbf{v}, \partial^* \phi)_{\mathbf{0}} &= (\mathbf{v}, \partial^\circ \phi)_{\mathbf{0}} + (\mathbf{v}, \mathbf{C}\phi)_{\mathbf{0}} \\ &= -\langle \partial^\circ \mathbf{v}, \phi \rangle + (\mathbf{v}, \mathbf{C}\phi)_{\mathbf{0}} - (\mathbf{v}, w_N \kappa \phi)_{\mathbf{0}} \\ &= -\langle \partial^* \mathbf{v}, \phi \rangle + (\mathbf{C}\mathbf{v}, \phi)_{\mathbf{0}} + (\mathbf{v}, \mathbf{C}\phi)_{\mathbf{0}} - (\mathbf{v}, w_N \kappa \phi)_{\mathbf{0}}, \end{aligned}$$

with $\mathbf{C} := \mathbf{AP}(\partial^{\circ} \mathbf{A}^{-1})$, and thus

$$\langle \partial_t(\phi_{-(\cdot)}\mathbf{v}), \widetilde{\boldsymbol{\xi}} \rangle = -\langle \partial^* \mathbf{v}, \mathbf{T} \boldsymbol{\xi} \rangle - (\mathbf{v}, w_N \kappa \, \mathbf{T} \boldsymbol{\xi} - (\mathbf{C} + \mathbf{C}^T) \mathbf{T} \boldsymbol{\xi} + (\partial^* \mathbf{T}) \boldsymbol{\xi})_{\mathbf{0}}.$$
(3.27)

Note that $\mathbf{T}: T\Gamma(t) \to T\Gamma(t)$ for all $t \in [0, T]$ and from Lemma 3.1 it follows that $\mathbf{T}, \mathbf{T}^{-1} \in C^1(\bar{S})^{3\times 3}$. Hence it holds

$$\mathbf{T}\boldsymbol{\xi} \in L^2_{H^1}$$
 and $\|\mathbf{T}\boldsymbol{\xi}\|_{\mathbf{1}} \simeq \|\boldsymbol{\xi}\|_{\mathbf{1}} \simeq \|\widetilde{\boldsymbol{\xi}}\|_{L^2([0,T],V_1(0))}$

From this, equality (3.27) and (3.24) one obtains after simple calculations

$$\begin{aligned} \|\partial_t(\phi_{-(.)}\mathbf{v})\|_{L^2([0,T],V_1(0)')} &\leq C(\|\partial^*\mathbf{v}\|_{L^2_{H^{-1}}} + \|\mathbf{v}\|_{\mathbf{0}}) \leq C(\|\partial^*\mathbf{v}\|_{L^2_{V'_1}} + \|\mathbf{v}\|_{\mathbf{0}}) \\ &\leq C\|\mathbf{v}\|_{W_*}. \end{aligned}$$

The reverse estimate $\|\partial^* \mathbf{v}\|_{L^2_{V'}} \leq C \|\mathbf{v}\|_{\mathcal{W}(V_1(0),V_1(0)')}$ follows from the identity

$$\langle \partial^* \mathbf{v}, \boldsymbol{\xi} \rangle = -\langle \partial_t(\phi_{-(\cdot)} \mathbf{v}), \phi_{-(\cdot)} \mathbf{T}^{-1} \boldsymbol{\xi} \rangle - (\mathbf{v}, w_N \kappa \, \boldsymbol{\xi} - (\mathbf{C} + \mathbf{C}^T) \boldsymbol{\xi} + (\partial^* \mathbf{T}) \mathbf{T}^{-1} \boldsymbol{\xi})_{\mathbf{0}}$$

by similar arguments (in particular the analogue result to (3.24) holds for the time derivative ∂_t on \mathcal{S}_0). Therefore we proved $\|\mathbf{v}\|_W \simeq \|\phi_{-(.)}\mathbf{v}\|_{\mathcal{W}(V_1(0),V_1(0)')}$ and hence $\mathbf{W}(V_1,V_1')$ and $\mathcal{W}(V_1(0),V_1(0)')$ are homeomorphic.

4. Well-Posed Weak Formulation

In this section, we introduce and analyze a weak formulation of TSNSE (2.12). We restrict our arguments to the solenoidal case g = 0. The extension of the analysis to the case $g \neq 0$ is discussed in Sec. 4.3. In the weak formulation we take a solution space with only solenoidal vector fields, and thus the pressure term vanishes. The existence of a corresponding unique pressure solution is shown in Sec. 4.1.

We introduce the notation

$$a(\mathbf{u}, \mathbf{v}) := 2\mu(E_s(\mathbf{u}), E_s(\mathbf{v}))_{\mathbf{0}}, \quad c(\mathbf{u}, \tilde{\mathbf{u}}, \mathbf{v}) := ((\nabla_{\Gamma} \mathbf{u})\tilde{\mathbf{u}}, \mathbf{v})_{\mathbf{0}},$$

$$\ell(\mathbf{u}, \mathbf{v}) := (w_N \mathbf{H} \mathbf{u}, \mathbf{v})_{\mathbf{0}},$$
(4.1)

and consider the following weak formulation of TSNSE (2.12) with g = 0: For given $\mathbf{f} \in L^2(\mathcal{S})^3$, with $\mathbf{f} = \mathbf{f}_T$, $\mathbf{u}_0 \in V_0(0)$, find $\mathbf{u}_T \in \mathbf{W}(V_1, V_1')$ such that $\mathbf{u}_T(0) = \mathbf{u}_0$ and

$$\langle \partial^{\circ} \mathbf{u}_T, \mathbf{v} \rangle + a(\mathbf{u}_T, \mathbf{v}) + c(\mathbf{u}_T, \mathbf{u}_T, \mathbf{v}) + \ell(\mathbf{u}_T, \mathbf{v}) = (\mathbf{f}, \mathbf{v})_{\mathbf{0}} \text{ for all } \mathbf{v} \in L^2_{V_1}.$$
 (4.2)

One easily checks that any smooth solution of (2.12) satisfies (4.2).

For the analysis of the weak formulation (4.2) we apply an established approach, e.g. see Ref. 43. Compared to the analysis of the non-stationary Navier–Stokes equations in Euclidean domains the main differences are that we use evolving spaces as introduced above instead of the standard Bochner ones, we have a normal time derivative ∂° in place of the usual $\frac{d}{dt}$, and an additional curvature-dependent term $(w_N \mathbf{H} \mathbf{u}_T, \mathbf{v})_0$. We show the existence of a Galerkin solution, derive *a priori* bounds and based on this show existence of a solution \mathbf{u}_T . We then show uniqueness of the solution with the help of Ladyzhenskaya's inequality.

Faedo-Galerkin approximation

The space $V_1(0)$ has a countable basis ψ_1, ψ_2, \ldots , which is pushed forward to a countable basis $\{\tilde{\psi}_i\}$ of $V_1(t)$ by letting $\tilde{\psi}_i = \phi_t \psi_i$. Consider

$$\mathbf{u}_m := \sum_{i=1}^m g_{i,m}(t) \widetilde{\boldsymbol{\psi}}_i. \tag{4.3}$$

We determine the unknown functions $g_{i,m}$ from (4.3) by considering the system of ODEs

$$(\partial^{\circ} \mathbf{u}_{m}, \widetilde{\boldsymbol{\psi}}_{j})_{0,t} + 2\mu(E_{s}(\mathbf{u}_{m}), E_{s}(\widetilde{\boldsymbol{\psi}}_{j}))_{0,t} + ((\nabla_{\Gamma}\mathbf{u}_{m})\mathbf{u}_{m}, \widetilde{\boldsymbol{\psi}}_{j})_{0,t} + (w_{N}\mathbf{H}\mathbf{u}_{m}, \widetilde{\boldsymbol{\psi}}_{j})_{0,t} = (\mathbf{f}, \widetilde{\boldsymbol{\psi}}_{j})_{0,t} \quad \text{for all } 1 \leq j \leq m.$$

$$\mathbf{u}_{m}(0) = \mathbf{u}_{0m}.$$

$$(4.4)$$

Here \mathbf{u}_{0m} is the L^2 -orthogonal projection of \mathbf{u}_0 on span{ ψ_1, \ldots, ψ_m }.

A priori bounds

Assume \mathbf{u}_m as in (4.3) satisfies (4.4). Multiplying (4.4) by $g_{j,m}(t)$ and summing over $j = 1, \ldots m$, we get, using $((\nabla_{\Gamma} \mathbf{u}_m) \mathbf{u}_m, \mathbf{u}_m)_{0,t} = 0$,

$$(\partial^{\circ} \mathbf{u}_{m}, \mathbf{u}_{m})_{0,t} + 2\mu(E_{s}(\mathbf{u}_{m}), E_{s}(\mathbf{u}_{m}))_{0,t} + (w_{N}\mathbf{H}\mathbf{u}_{m}, \mathbf{u}_{m})_{0,t} = (\mathbf{f}, \mathbf{u}_{m})_{0,t}, \quad (4.5)$$

and applying integration by parts (3.17), we have

$$\frac{d}{dt} \|\mathbf{u}_m\|_{0,t}^2 + 4\mu \|E_s(\mathbf{u}_m)\|_{0,t}^2 = -2(w_N \mathbf{H} \mathbf{u}_m, \mathbf{u}_m)_{0,t} + (w_N \kappa \mathbf{u}_m, \mathbf{u}_m)_{0,t} + 2(\mathbf{f}, \mathbf{u}_m)_{0,t}.$$
(4.6)

From this we obtain for $0 < \tau \leq T$,

$$\|\mathbf{u}_{m}\|_{0,\tau}^{2} + 4\mu \int_{0}^{\tau} \|E_{s}(\mathbf{u}_{m})\|_{0,t}^{2} dt \lesssim \int_{0}^{\tau} \|\mathbf{u}_{m}\|_{0,t}^{2} dt + \int_{0}^{\tau} \|\mathbf{f}\|_{0,t}^{2} dt + \|\mathbf{u}_{0m}\|_{L^{2}(\Gamma_{0})}^{2}.$$
(4.7)

Here and in the remainder we write $A \leq B$ to denote $A \leq c B$ with some constant c which may depend on the final time T, the maximum normal velocity $||w_N||_{L^{\infty}(S)}$ and on smoothness properties of the space-time manifold, quantified by $||\mathbf{H}||_{L^{\infty}(S)}$.

Note that $\|\kappa\|_{L^{\infty}(\mathcal{S})} = \|\operatorname{tr}(\mathbf{H})\|_{L^{\infty}(\mathcal{S})} \leq 2\|\mathbf{H}\|_{L^{\infty}(\mathcal{S})}$. The Gronwall lemma and (4.7) yield the *a priori* bound,

$$\max_{0 \le t \le T} \|\mathbf{u}_m\|_{0,t} + \|E_s(\mathbf{u}_m)\|_{\mathbf{0}} \lesssim \|\mathbf{f}\|_{\mathbf{0}} + \|\mathbf{u}_0\|_{L^2(\Gamma_0)}.$$
(4.8)

The uniform Korn inequality and the estimates in (4.7) and (4.8) yield the *a priori* estimate

$$\|\mathbf{u}_m\|_{\mathbf{1}} \lesssim \|\mathbf{f}\|_{\mathbf{0}} + \|\mathbf{u}_0\|_{L^2(\Gamma_0)}.$$
(4.9)

Existence of solution

Consider the ODEs system (4.4). Due to Lemma 3.6 we have $\mathbf{P}\partial^{\circ}(\tilde{\psi}_i) = \partial^*(\tilde{\psi}_i) - \mathbf{C}\tilde{\psi}_i = -\mathbf{C}\tilde{\psi}_i$, with $\mathbf{C} = \mathbf{A}\mathbf{P}\partial^{\circ}\mathbf{A}^{-1}$. Thus (4.4) results in the following system for $g_{i,m}$:

$$\sum_{i=1}^{m} \frac{dg_{i,m}(t)}{dt} (\widetilde{\psi}_{i}, \widetilde{\psi}_{j})_{0,t} = -\sum_{i,k=1}^{m} g_{i,m}(t)g_{k,m}(t)((\nabla_{\Gamma}\widetilde{\psi}_{k})\widetilde{\psi}_{i}, \widetilde{\psi}_{j})_{0,t}$$
$$-\sum_{i=1}^{m} g_{i,m}(t)\{2\mu(E_{s}(\widetilde{\psi}_{i}), E_{s}(\widetilde{\psi}_{j}))_{0,t} + ((w_{N}\mathbf{H} - \mathbf{C})\widetilde{\psi}_{i}, \widetilde{\psi}_{j})_{0,t}\} + (\mathbf{f}, \widetilde{\psi}_{j})_{0,t},$$

$$(4.10)$$

for $1 \leq j \leq m$. From the fact that the pushforward map ϕ_t is one-to-one and linear for every t, and ψ_i are linear independent we infer that $\tilde{\psi}_i$ are linear independent for every t and thus the matrix $\mathbf{M}(t) := (\tilde{\psi}_i(t), \tilde{\psi}_j(t))_{1 \leq i,j \leq m}$ is invertible for $t \in [0, T]$. Moreover, (3.6) and the definition of $\tilde{\psi}$ implies $\mathbf{M} \in C^1[0, T]^{m \times m}$. Since any eigenvalue of \mathbf{M} , denoted by $\lambda(\mathbf{M})$, continuously depends on matrix coefficients, the bound $\lambda(\mathbf{M}) > 0$ for each $t \in [0, T]$ implies $\lambda(\mathbf{M}) \geq c > 0$ uniformly on [0, T]. The uniform lower bound for the eigenvalues and the symmetry of \mathbf{M} ensures $\|\mathbf{M}^{-1}\|_{C[0,T]} \leq C$. Multiplying both sides of (4.10) with \mathbf{M}^{-1} , one verifies that the Picard–Lindelöf theorem applies. Hence a unique solution $g_{i,m}(t), 1 \leq i \leq m$, exists for a maximal interval $[0, t_e], t_e > 0$. If $t_e < T$, then $\lim_{t \uparrow t_e} \|\mathbf{u}_m(t)\|_{0,t} = \infty$, which contradicts the established bound (4.8) with T replaced by t_e . Hence, a unique solution $\mathbf{u}_m(t)$ exists for $t \in [0, T]$.

From the *a priori* bounds (4.8), (4.9) it follows that there is a subsequence of $(\mathbf{u}_{m'})_{m'\geq 1}$ of $(\mathbf{u}_m)_{m\geq 1}$ that is weak-star convergent in L^{∞}_H and weakly convergent in $L^2_{V_1}$ to $\mathbf{u}^* \in L^{\infty}_H \cap L^2_{V_1}$. Due to the compactness of $L^2_{V_1} \hookrightarrow L^2_{V_0}$ this sequence also strongly converges in $L^2_{V_0}$. Now note that, with $\tilde{\psi}_j$, $j = 1, 2, \ldots$, as above, functions $\sum_{j=1}^N g_j(t) \tilde{\psi}_j$, $N \in \mathbb{N}$, $g_j \in C^1([0,T])$, with $g_j(T) = 0$, are dense in $L^2_{V_1}$. We multiply (4.4) with such a function g_j , integrate over [0,T] and apply partial integration (3.17), which yields

$$-(\mathbf{u}_{m'},\partial^{\circ}(\widehat{\boldsymbol{\psi}}_{j}g_{j}))_{\mathbf{0}} - (\kappa w_{N}\mathbf{u}_{m'},\widehat{\boldsymbol{\psi}}_{j}g_{j})_{\mathbf{0}}$$

$$\tag{4.11}$$

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$$+ 2\mu(E_s(\mathbf{u}_{m'}), E_s(\widetilde{\psi}_j g_j))_{\mathbf{0}} + (\nabla_{\Gamma} \mathbf{u}_{m'} \mathbf{u}_{m'}, \widetilde{\psi}_j g_j)_{\mathbf{0}}$$

$$(4.12)$$

$$+(w_N \mathbf{H} \mathbf{u}_{m'}, \widetilde{\psi}_j g_j)_{\mathbf{0}} = (\mathbf{u}_{0m}, \psi_j g_j(0))_{L^2(\Gamma_0)} + (\mathbf{f}, \widetilde{\psi}_j g_j)_{\mathbf{0}}.$$
(4.13)

Due to the strong convergence $\mathbf{u}_{m'} \to \mathbf{u}^*$ in $L^2_{V_0}$ we can pass to the limit in the two terms in (4.11) and the first term in (4.13). Since $(E_s(\cdot), E_s(\mathbf{v}))_0$ is a functional on $L^2_{V_1}$ for any $\mathbf{v} \in L^2_{V_1}$, we can pass to the limit in the first term in (4.12). Using the strong convergence in $L^2_{V_0}$ we can also pass to the limit in the second term in (4.12), cf. Ref. 43. By definition of \mathbf{u}_{0m} we have $\mathbf{u}_{0m} \to \mathbf{u}(0)$ strongly in $L^2(\Gamma_0)$. Thus we get, cf. (4.1),

$$-(\mathbf{u}^*, \partial^{\circ}(\widetilde{\boldsymbol{\psi}}_j g_j))_{\mathbf{0}} - (\kappa w_N \mathbf{u}^*, \widetilde{\boldsymbol{\psi}}_j g_j)_{\mathbf{0}} = -a(\mathbf{u}^*, \widetilde{\boldsymbol{\psi}}_j g_j) - c(\mathbf{u}^*, \mathbf{u}^*, \widetilde{\boldsymbol{\psi}}_j g_j) - \ell(\mathbf{u}^*, \widetilde{\boldsymbol{\psi}}_j g_j) + (\mathbf{u}_0, \boldsymbol{\psi}_j g_j(0))_{L^2(\Gamma_0)} + (\mathbf{f}, \widetilde{\boldsymbol{\psi}}_j g_j)_{\mathbf{0}}.$$
(4.14)

We restrict to g_i with $g_i(0) = 0$ and build linear combinations of (4.14) to arrive at

$$-(\mathbf{u}^*, \partial^\circ \mathbf{v})_{\mathbf{0}} - (\kappa w_N \mathbf{u}^*, \mathbf{v})_{\mathbf{0}} = -a(\mathbf{u}^*, \mathbf{v}) - c(\mathbf{u}^*, \mathbf{u}^*, \mathbf{v}) - \ell(\mathbf{u}^*, \mathbf{v}) + (\mathbf{f}, \mathbf{v})_{\mathbf{0}}$$

$$(4.15)$$

for all $\mathbf{v} = \sum_{j=1}^{N} \widetilde{\psi}_j g_j$. We estimate the nonlinear term with the help of uniform Ladyzhenskaya inequality and (4.8), (4.9):

$$c(\mathbf{u}^{*}, \mathbf{u}^{*}, \mathbf{v}) = -c(\mathbf{v}, \mathbf{u}^{*}, \mathbf{u}^{*}) \leq \int_{0}^{T} \|\nabla_{\Gamma} \mathbf{v}\|_{0,t} \|\mathbf{u}^{*}\|_{L^{4}(\Gamma(t))}^{2} dt$$

$$\lesssim \int_{0}^{T} \|\mathbf{v}\|_{1,t} \|\mathbf{u}^{*}\|_{0,t} \|\mathbf{u}^{*}\|_{1,t} dt \lesssim \sup_{t \in [0,T]} \|\mathbf{u}^{*}\|_{0,t} \|\mathbf{u}^{*}\|_{1} \|\mathbf{v}\|_{1}$$

$$\lesssim (\|\mathbf{f}\|_{0} + \|\mathbf{u}_{0}\|_{L^{2}(\Gamma_{0})})^{2} \|\mathbf{v}\|_{1}.$$
(4.16)

Using the above estimate and obvious continuity estimates for other terms on the right-hand side in (4.15) together with a density argument we conclude that

$$\|\partial^{\circ} \mathbf{u}^{*}\|_{L^{2}_{V'_{1}}} \lesssim F(1+F), \text{ with } F := \|\mathbf{f}\|_{\mathbf{0}} + \|\mathbf{u}_{0}\|_{L^{2}(\Gamma_{0})},$$
 (4.17)

hence $\mathbf{u}^* \in \mathbf{W}(V_1, V_1')$ and furthermore $\mathbf{u}_T = \mathbf{u}^*$ satisfies (4.2).

To check that $\mathbf{u}^*(0) = \mathbf{u}_0$ holds, we apply standard arguments. Using continuity of $t \to \|\mathbf{v}(t)\|_{0,t}$ for $\mathbf{v} \in \mathbf{W}(V_1, V_1')$ (1) and density of smooth functions in $\mathbf{W}(V_1, V_1')$ it follows that the partial integration rule (3.17) can be generalized to $\mathbf{W}(V_1, V_1')$. Test (4.2) with $\mathbf{v} = \tilde{\psi}_j g_j(t)$, with $g_j(0) = 1$, applying partial integration and comparing the result with (4.14) we obtain $(\mathbf{u}^*, \psi_j)_{L^2(\Gamma_0)} = (\mathbf{u}_0, \psi_j)_{L^2(\Gamma_0)}$. Since $(\psi_j)_{j\in\mathbb{N}}$ is dense in $V_0(0)$ we conclude that $\mathbf{u}^*(0) = \mathbf{u}_0$ holds.

Uniqueness of solution

We prove uniqueness of the solution using essentially the same arguments as in Euclidean space. For the sake of presentation below, we use $\langle \langle \cdot, \cdot \rangle \rangle$ to denote L^2

duality pairing between $V_1(t)$ and $V_1(t)'$ and introduce the notation for t-level bilinear forms, cf. (4.1):

$$a_t(\mathbf{v}, \boldsymbol{\psi}) := 2\mu(E_s(\mathbf{v}), E_s(\boldsymbol{\psi}))_{0,t}, \quad c_t(\mathbf{v}, \tilde{\mathbf{v}}, \boldsymbol{\psi}) := ((\nabla_{\Gamma} \mathbf{v})\tilde{\mathbf{v}}, \boldsymbol{\psi})_{0,t},$$
$$\ell_t(\mathbf{v}, \boldsymbol{\psi}) := (w_N \mathbf{H} \mathbf{v}, \boldsymbol{\psi})_{0,t}.$$

Note that $c_t(\mathbf{v}, \tilde{\mathbf{v}}, \mathbf{v}) = 0$ holds. A solution $\mathbf{u}_T \in \mathbf{W}(V_1, V_1')$ of (4.2) satisfies

$$\langle\!\langle \partial^{\circ} \mathbf{u}_{T}(t), \mathbf{v}(t) \rangle\!\rangle + a_{t}(\mathbf{u}_{T}(t), \mathbf{v}(t)) + c_{t}(\mathbf{u}_{T}(t), \mathbf{u}_{T}(t), \mathbf{v}(t)) + \ell_{t}(\mathbf{u}_{T}(t), \mathbf{v}(t)) = (\mathbf{f}(t), \mathbf{v}(t))_{0,t} \quad \text{a.e. in } t \quad \text{for } \mathbf{v} \in L^{2}_{V_{1}}.$$

$$(4.18)$$

The Leibniz rule (3.16) extends to $\mathbf{v} \in \mathbf{W}(V_1, V_1')$ (cf. Ref. 1) yielding

$$\frac{d}{dt}(\mathbf{v}(t),\mathbf{v}(t))_{0,t} = 2\langle\!\langle \partial^{\circ} \mathbf{v}(t),\mathbf{v}(t) \rangle\!\rangle + \int_{\Gamma(t)} |\mathbf{v}(t)|^2 w_N \kappa \, ds.$$

Let \mathbf{u}_T^1 , \mathbf{u}_T^2 be solutions of (4.2) with $\mathbf{u}_T^1(0) = \mathbf{u}_T^2(0) = \mathbf{u}_0$. Letting $\boldsymbol{\psi} := \mathbf{u}_T^1 - \mathbf{u}_T^2$ and using $c_t(\boldsymbol{\psi}, \mathbf{u}_T^1, \boldsymbol{\psi}) = 0$ we compute, with $C_1 := \|w_N \kappa\|_{C(\mathcal{S})}$,

$$\begin{split} \frac{d}{dt} \|\psi(t)\|_{0,t}^{2} + 2a_{t}(\psi(t),\psi(t)) \\ &= 2\langle\!\langle\partial^{\circ}\psi(t),\psi(t)\rangle\!\rangle + \int_{\Gamma(t)} |\psi(t)|^{2}w_{N}\kappa\,ds + 2a_{t}(\psi(t),\psi(t)) \\ &\leq 2\langle\!\langle\partial^{\circ}\mathbf{u}_{T}^{1}(t),\psi(t)\rangle\!\rangle + 2a_{t}(\mathbf{u}_{T}^{1}(t),\psi(t)) \\ &- 2\langle\!\langle\partial^{\circ}\mathbf{u}_{T}^{2}(t),\psi(t)\rangle\!\rangle - 2a_{t}(\mathbf{u}_{T}^{2}(t),\psi(t)) + C_{1}\|\psi(t)\|_{0,t}^{2} \\ &= -2c_{t}(\mathbf{u}_{T}^{1}(t),\mathbf{u}_{T}^{1}(t),\psi(t)) - 2\ell_{t}(\mathbf{u}_{T}^{1}(t),\psi(t)) \\ &+ 2c_{t}(\mathbf{u}_{T}^{2}(t),\mathbf{u}_{T}^{2}(t),\psi(t)) - 2\ell_{t}(\mathbf{u}_{T}^{2}(t),\psi(t)) + C_{1}\|\psi(t)\|_{0,t}^{2} \\ &= -2c_{t}(\mathbf{u}_{T}^{2}(t),\mathbf{u}_{T}(t),\psi(t)) - 2\ell_{t}(\mathbf{u}_{T}(t),\psi(t)) + C_{1}\|\psi(t)\|_{0,t}^{2} \end{split}$$

We have $|\ell_t(\boldsymbol{\psi}(t), \boldsymbol{\psi}(t))| \leq ||w_N \mathbf{H}||_{C(S)} ||\boldsymbol{\psi}(t)||_{0,t}^2$. For the other terms on the righthand side above we use (3.14) and the Korn inequality (3.11) to estimate

$$\begin{split} &\frac{d}{dt} \|\boldsymbol{\psi}(t)\|_{0,t}^{2} + 2a_{t}(\boldsymbol{\psi}(t),\boldsymbol{\psi}(t)) \\ &\leq C \|\mathbf{u}_{T}^{2}\|_{1,t} \|\boldsymbol{\psi}(t)\|_{L^{4}(\Gamma(t))}^{2} + C \|\boldsymbol{\psi}(t)\|_{0,t}^{2} \\ &\leq C \|\mathbf{u}_{T}^{2}\|_{1,t} \|\boldsymbol{\psi}(t)\|_{0,t} \|\boldsymbol{\psi}\|_{1,t} + C \|\boldsymbol{\psi}(t)\|_{0,t}^{2} \\ &\leq C \|\mathbf{u}_{T}^{2}\|_{1,t} \|\boldsymbol{\psi}(t)\|_{0,t} (\|\boldsymbol{\psi}(t)\|_{0,t} + a_{t}(\boldsymbol{\psi}(t),\boldsymbol{\psi}(t))^{\frac{1}{2}}) + C \|\boldsymbol{\psi}(t)\|_{0,t}^{2} \\ &\leq C (\|\mathbf{u}_{T}^{2}\|_{1,t} + \|\mathbf{u}_{T}^{2}\|_{1,t}^{2}) \|\boldsymbol{\psi}(t)\|_{0,t}^{2} + 2a_{t}(\boldsymbol{\psi}(t),\boldsymbol{\psi}(t)), \end{split}$$

with a suitable constant C independent of $t \in [0, T]$ and of $\mathbf{u}_T^1, \mathbf{u}_T^2$. Thus, we get

$$\frac{d}{dt}\|\boldsymbol{\psi}(t)\|_{0,t}^2 \le f_u(t)\|\boldsymbol{\psi}(t)\|_{0,t}^2, \quad f_u(t) := C(\|\mathbf{u}_T^2\|_{1,t} + \|\mathbf{u}_T^2\|_{1,t}^2).$$

Now, $\mathbf{u}_T^2 \in L^2_{V_1}$ implies that $\int_0^T f_u(s) ds$ is bounded and so the Gronwall inequality together with $\|\psi(0)\|_{L^2(\Gamma_0)} = 0$ yields $\psi(t) = 0$ for $t \in [0, T]$ and thus the uniqueness result holds.

Summarizing we proved the following main well-posedness result.

Theorem 4.1. The weak formulation (4.2) of the TSNSE has a unique solution $\mathbf{u}_T \in \mathbf{W}(V_1, V_1')$. The solution satisfies

$$\|\mathbf{u}_T\|_W \le C(1+F)F, \quad \text{with } F := \|\mathbf{f}\|_{\mathbf{0}} + \|\mathbf{u}_0\|_{L^2(\Gamma_0)}.$$
 (4.19)

4.1. Surface pressure

For $\mathbf{v} \in L^2_{V_1}$, (3.18) defines $\partial^{\circ} \mathbf{v}$ as a functional on $\widehat{\mathcal{D}}_0$. The density of $\widehat{\mathcal{D}}_0$ in $L^2_{H^1}$ and the density of $\mathcal{D}_0 \subset \widehat{\mathcal{D}}_0$ in $L^2_{V_1}$ is used to define the bounded linear functionals $\partial^{\circ} \mathbf{v} \in L^2_{H^{-1}}$ and $\partial^{\circ} \mathbf{v} \in L^2_{V'}$, respectively. The following equivalence holds:

$$\partial^{\circ} \mathbf{v} \in L^2_{V'_1} \Leftrightarrow \partial^{\circ} \mathbf{v} \in L^2_{H^{-1}}, \quad \mathbf{v} \in L^2_{V_1}.$$

$$(4.20)$$

Implication " \Leftarrow " in (4.20) is trivial since $V_1 \subset H^1$. The " \Rightarrow " implication follows from (3.24) and Lemma 3.6.

Below we introduce a weak formulation of TSNSE on the velocity space

$$\begin{split} \mathbf{W}(H^1, H^{-1}) &= \{ \mathbf{v} \in L^2_{H^1} \, : \, \partial^{\circ} \mathbf{v} \in L^2_{H^{-1}} \} \\ & \text{with } (\cdot, \cdot)_{W(H^1, H^{-1})} = (\cdot, \cdot)_{\mathbf{1}} + (\cdot, \cdot)_{L^2_{H^{-1}}}, \end{split}$$

with a pressure unknown $\pi \in L^2(\mathcal{S})$. One checks that $\mathbf{W}(H^1, H^{-1})$ is a Hilbert space by the same arguments as for $\mathbf{W}(V_1, V_1')$. Consider the following mixed formulation of TSNSE, which relates to the well-posed weak formulation (4.2): For given $\mathbf{f} \in L^2(\mathcal{S})^3$, with $\mathbf{f} = \mathbf{f}_T$, $\mathbf{u}_0 \in V_0(0)$, find $\mathbf{u}_T \in \mathbf{W}(H^1, H^{-1})$ and $\pi \in L^2(\mathcal{S})$, with $\int_{\Gamma(t)} \pi \, ds = 0$ a.e. $t \in [0, T]$, such that $\mathbf{u}_T(0) = \mathbf{u}_0$ and

$$\langle \partial^{\circ} \mathbf{u}_{T}, \mathbf{v} \rangle + a(\mathbf{u}_{T}, \mathbf{v}) + c(\mathbf{u}_{T}, \mathbf{u}_{T}, \mathbf{v}) + \ell(\mathbf{u}_{T}, \mathbf{v}) + (\pi, \operatorname{div}_{\Gamma} \mathbf{v})_{\mathbf{0}} = (\mathbf{f}, \mathbf{v})_{\mathbf{0}}$$

$$(q, \operatorname{div}_{\Gamma} \mathbf{u}_{T})_{\mathbf{0}} = 0$$

$$(4.21)$$

for all $\mathbf{v} \in L^2_{H^1}$, $q \in L^2(\mathcal{S})$.

Theorem 4.2. The problem (4.21) has a unique solution (\mathbf{u}_T, π) . The velocity solution \mathbf{u}_T is also the unique solution of (4.2). Furthermore, with $F := \|\mathbf{f}\|_{\mathbf{0}} + \|\mathbf{u}_0\|_{L^2(\Gamma_0)}$ the following estimate holds:

$$\|\mathbf{u}_T\|_W + \|\pi\|_{\mathbf{0}} \le C(1+F)F.$$
(4.22)

Proof. Let $\mathbf{u}_T \in \mathbf{W}(V_1, V_1')$ be the solution of (4.2). Define

$$\widetilde{\mathbf{f}}(\mathbf{v}) := \langle \partial^{\circ} \mathbf{u}_{T}, \mathbf{v} \rangle + a(\mathbf{u}_{T}, \mathbf{v}) + c(\mathbf{u}_{T}, \mathbf{u}_{T}, \mathbf{v}) + \ell(\mathbf{u}_{T}, \mathbf{v}) - (\mathbf{f}, \mathbf{v})_{\mathbf{0}}, \quad \mathbf{v} \in L^{2}_{H^{1}}.$$

Using (4.20) and straightforward estimates we obtain $\tilde{\mathbf{f}} \in L^2_{H^{-1}}$.^b We use the standard argument (e.g. Remark I.1.9 in Ref. 43) that for every $t \in [0,T]$ estimate (3.13) implies that $\nabla_{\Gamma} \in \mathcal{L}(L^2(\Gamma(t)), H^{-1}(t))$ has a closed range $R(\nabla_{\Gamma})$ in $H^{-1}(t)$ and so

$$R(\nabla_{\Gamma}) = \ker(\nabla_{\Gamma}^*)^{\perp}$$
, with $\ker(\nabla_{\Gamma}^*) = V_1(t)$.

Note that $\widetilde{\mathbf{f}}(t)$ is an element of $H^{-1}(t)$ for a.e. $t \in [0, T]$ and, since \mathbf{u}_T is the solution of (4.2), $\langle \langle \widetilde{\mathbf{f}}(t), \mathbf{v} \rangle = 0$ for all $\mathbf{v} \in V_1(t)$. Hence, $\widetilde{\mathbf{f}}(t) \in R(\nabla_{\Gamma})$ which means

 $\widetilde{\mathbf{f}}(t) = \nabla_{\Gamma} \pi(t) \text{ for some } \pi(t) \in L^2(\Gamma(t)) \text{ for a.e. } t \in [0, T].$

We take $\pi(t)$ such that $\int_{\Gamma(t)} \pi(t) = 0$ holds. To see that $t \to ||\pi(t)||_{0,t}$ is measurable we argue as follows. First note that $t \to ||\pi(t)||_{0,t}$ is measurable if $t \to \pi_0(t)$ is Bochner-measurable on [0,T] for $\pi_0 := J\pi \circ \Phi_t^n$. Since $L^2(\Gamma_0)$ is separable, it suffices to check that $t \to (\pi_0(t), q)_{0,0}$ is a measurable function for any $q \in L^2(\Gamma_0)$ such that $\int_{\Gamma_0} q = 0$. Property (3.13) for t = 0 implies that there exists $\mathbf{v}_0 \in H^1(0)$ such that $q = \operatorname{div}_{\Gamma} \mathbf{v}_0$ and $||\mathbf{v}_0||_{1,0} \leq C_0 ||q||_{0,0}$. Let $\mathbf{v} = \phi_t \mathbf{v}_0$. We then have

$$(\pi_0(t), q)_{0,0} = (\pi_0(t), \operatorname{div}_{\Gamma} \mathbf{v}_0)_{0,0} = (\pi(t), \operatorname{div}_{\Gamma} \mathbf{v})_{0,t} = \langle \nabla_{\Gamma} \pi(t), \mathbf{v} \rangle_{H^{-1}(t) \times H^1(t)},$$

which implies that $t \to \langle \nabla_{\Gamma} \pi(t), \mathbf{v} \rangle_{H^{-1}(t) \times H^{1}(t)}$ is measurable, since $\nabla_{\Gamma} \pi \in L^{2}_{H^{-1}}$ and $\mathbf{v} \in L^{2}_{H^{1}}$.

From (4.21) we get, with notation as in (4.1),

$$(\pi, \operatorname{div}_{\Gamma} \mathbf{v})_{0,t} = (\mathbf{f}(t), \mathbf{v}(t))_{0,t} - \langle\!\langle \partial^{\circ} \mathbf{u}_{T}(t), \mathbf{v}(t) \rangle\!\rangle - a_{t}(\mathbf{u}_{T}(t), \mathbf{v}(t)) - c_{t}(\mathbf{u}_{T}(t), \mathbf{u}_{T}(t), \mathbf{v}(t)) - \ell_{t}(\mathbf{u}_{T}(t), \mathbf{v}(t)) \quad \text{a.e. in } t \text{ for } \mathbf{v} \in L^{2}_{V_{1}}.$$

Using the uniform inf-sup estimate, cf. Lemma 3.3, we get

$$\|\pi(t)\|_{0,t} \le C(\|\mathbf{f}(t)\|_{0,t} + \|\partial^{\circ}\mathbf{u}_{T}\|_{V_{1}(t)'} + \|\mathbf{u}_{T}(t)\|_{1,t}(1 + \|\mathbf{u}_{T}(t)\|_{0,t}))$$

with a constant C independent of t. Hence, $\pi \in L^2(\mathcal{S})$ holds. The estimate for velocity in (4.22) is the same as in Theorem 4.19. Note that $\max_{t \in [0,T]} \|\mathbf{u}_T(t)\|_{0,t} \lesssim$ F holds, cf. (4.8). Using this and the velocity estimate we obtain the bound for the pressure in (4.22). Uniqueness of \mathbf{u}_T follows by restricting to $\mathbf{v} \in L^2_{V_1}$ in (4.21) and using the fact that (4.2) has a unique solution. Uniqueness of π is easily derived using the inf-sup property.

4.2. Energy balance

Multiplying (2.12) by \mathbf{u}_T , integrating over $\Gamma(t)$ and using (3.16), we obtain, for a smooth solution, the energy balance of the system at any $t \in (0, T)$,

$$\frac{1}{2}\frac{d}{dt}\|\mathbf{u}_T\|_{0,t}^2 + 2\mu\|E_s(\mathbf{u}_T)\|_{0,t}^2 + \left(w_N\left(\mathbf{H} - \frac{1}{2}\kappa\mathbf{I}\right)\mathbf{u}_T, \mathbf{u}_T\right)_{0,t} = (\mathbf{f}, \mathbf{u}_T)_{0,t}.$$
(4.23)

^bTo see $\mathbf{v} \to c(\mathbf{u}_T, \mathbf{u}_T, \mathbf{v}) \in L^2_{H^{-1}}$, one uses the same arguments as in (4.16).

Next, we comment on the contribution of the third term in (4.23), which appears if the surface is both evolving and non-flat. First, we note that $\mathbf{H} - \frac{1}{2}\kappa \mathbf{I} = \mathbf{H} - \frac{1}{2}\kappa \mathbf{P}$ on $T\Gamma(t)$ and $\mathbf{H} - \frac{1}{2}\kappa \mathbf{P} = \mathbf{H} - \frac{1}{2}\mathrm{tr}(\mathbf{H})\mathbf{P}$. Since $\mathrm{tr}(\mathbf{P}) = 2$ we get $\mathrm{tr}(\mathbf{H} - \frac{1}{2}\kappa \mathbf{P}) = 0$. This implies that the symmetric matrix $\mathbf{H} - \frac{1}{2}\kappa \mathbf{P}$ has real eigenvalues $\{0, \lambda, -\lambda\}$ with the 0 eigenvalue corresponding to vectors normal to $\Gamma(t)$. Take a fixed point \mathbf{x} on the surface $\Gamma(t)$ with $w_N(\mathbf{x}) \neq 0$. Denote by κ_1 and κ_2 the two principle curvatures of $\Gamma(t)$. For the eigenvalue of $\mathbf{H}(\mathbf{x}) - \frac{1}{2}\kappa(\mathbf{x})\mathbf{P}(\mathbf{x})$ we have $\lambda(\mathbf{x}) = \frac{1}{2}(\kappa_1(\mathbf{x}) - \kappa_2(\mathbf{x}))$. Therefore $w_N(\mathbf{x})(\mathbf{H}(\mathbf{x}) - \frac{1}{2}\kappa(\mathbf{x})\mathbf{P}(\mathbf{x})) = 0$ iff $\kappa_1(\mathbf{x}) = \kappa_2(\mathbf{x})$ holds, and it is indefinite otherwise. In the latter case the increase or decrease of kinetic energy due to this term depends on the alignment of the flow with the principle directions and the sign of w_N .

4.3. Non-solonoidal solution

The tangential surface Navier–Stokes system (2.12) admits non-solonoidal solutions with $\operatorname{div}_{\Gamma} \mathbf{u}_T = g$, where $g = -w_N \kappa$, $\int_{\Gamma(t)} g \, ds = 0$ for $t \in [0, T]$, is defined by the surface geometry and evolution. We outline how the analysis for the solonoidal case presented above can be extended to such a problem. We assume that $g : S \to \mathbb{R}$ is sufficiently regular. Let $\phi(t, \mathbf{x})$ be the unique solution of the Laplace–Beltrami equation $\Delta_{\Gamma(t)}\phi = g$, $\int_{\Gamma(t)} \phi ds = 0$, and define $\tilde{\mathbf{u}}_T := \nabla_{\Gamma}\phi$. Then (\mathbf{u}_T, π) solves (2.12) iff $\hat{\mathbf{u}}_T = \mathbf{u}_T - \tilde{\mathbf{u}}_T$ and π solve the system

$$\begin{cases} \mathbf{P}\partial^{\circ}\widehat{\mathbf{u}}_{T} + (\nabla_{\Gamma}\widehat{\mathbf{u}}_{T})\widehat{\mathbf{u}}_{T} + (\nabla_{\Gamma}\widetilde{\mathbf{u}}_{T})\widehat{\mathbf{u}}_{T} + (\nabla_{\Gamma}\widehat{\mathbf{u}}_{T})\widetilde{\mathbf{u}}_{T} \\ + w_{N}\mathbf{H}\widehat{\mathbf{u}}_{T} - 2\mu\mathbf{P}\operatorname{div}_{\Gamma}E_{s}(\widehat{\mathbf{u}}_{T}) + \nabla_{\Gamma}\pi = \widehat{\mathbf{f}} \\ \operatorname{div}_{\Gamma}\widehat{\mathbf{u}}_{T} = 0, \end{cases}$$
(4.24)

with

$$\widehat{\mathbf{f}} = \mathbf{f} - (\mathbf{P}\partial^{\circ}\widetilde{\mathbf{u}}_{T} + (\nabla_{\Gamma}\widetilde{\mathbf{u}}_{T})\widetilde{\mathbf{u}}_{T} + w_{N}\mathbf{H}\widetilde{\mathbf{u}}_{T} - 2\mu\mathbf{P}\operatorname{div}_{\Gamma}E_{s}(\widetilde{\mathbf{u}}_{T})).$$

The two additional terms $(\nabla_{\Gamma} \widetilde{\mathbf{u}}_T) \widehat{\mathbf{u}}_T$ and $(\nabla_{\Gamma} \widehat{\mathbf{u}}_T) \widetilde{\mathbf{u}}_T$ in the momentum equation in (4.24) are linear with respect to the unknown velocity field $\widehat{\mathbf{u}}_T$ and can be treated very similar to the zero order term $w_N \mathbf{H} \widehat{\mathbf{u}}_T$. The necessary regularity of $\widehat{\mathbf{f}}$ can be established using the smoothness of g and \mathcal{S} . We skip working out further details.

5. Discretization Method

As discussed in the introduction, only very few papers are available in which finite element discretization methods for vector- or tensor-valued surface PDEs, such as the surface Navier–Stokes equations, on evolving surfaces are treated. In this section, we present a discretization method for the TSNSE (2.12). The method is based on a combination of an implicit time stepping scheme with a TraceFEM for discretization in space. The general idea behind TraceFEM is to use standard time-independent (bulk) finite element spaces to approximate surface quantities. The method is based on tangential calculus in the ambient space $\mathbb{R}^3 \supset \Gamma(t)$. For scalar PDEs on evolving surfaces, a space-time variant of TraceFEM is treated in Ref. 30. A finite difference (FD) in time — FEM in space variant for PDEs on timedependent surfaces is treated in Ref. 23 (scalar problems) and in time-dependent volumetric domains in Ref. 22 (scalar equations) and Ref. 46 (Stokes problem). Compared to the space-time variant the FD-FEM approach is more flexible in terms of implementation and the choice of elements. Below we explain this FD-FEM approach applied to the TSNSE. We start with the numerical treatment of the system's evolution in time.

5.1. Time-stepping scheme

Consider uniformly distributed time nodes $t_n = n\Delta t$, n = 0, ..., N, with the time step $\Delta t = T/N$. We assume that the time step Δt is sufficiently small such that

$$\Gamma(t_n) \subset \mathcal{O}(\Gamma(t_{n-1})), \quad n = 1, \dots, N,$$
(5.1)

with $\mathcal{O}(\Gamma(t))$ a neighborhood of $\Gamma(t)$ where a smooth extension of surface quantities on $\Gamma(t)$ is well defined. Assuming a smooth extension \mathbf{u}_T^e , we rewrite the normal time derivative ∂° used in (2.12) in terms of standard time and space derivatives:

$$\mathbf{P}\partial^{\circ}\mathbf{u}_{T} + (\nabla_{\Gamma}\mathbf{u}_{T})\mathbf{u}_{T} = \mathbf{P}\left(\frac{\partial\mathbf{u}_{T}^{e}}{\partial t} + (\nabla\mathbf{u}_{T}^{e})\mathbf{w}_{N} + (\nabla\mathbf{u}_{T}^{e})\mathbf{u}_{T}\right)$$
$$= \mathbf{P}\left(\frac{\partial\mathbf{u}_{T}^{e}}{\partial t} + (\nabla\mathbf{u}_{T}^{e})(\mathbf{w}_{N} + \mathbf{u}_{T})\right).$$
(5.2)

On $\Gamma(t_n)$ the time derivative term is approximated by

$$\mathbf{P}\frac{\partial \mathbf{u}_T^e}{\partial t} \approx \frac{\mathbf{u}_T(t_n) - \mathbf{P}(t_n)\mathbf{u}_T(t_{n-1})^e}{\Delta t}.$$

Note that due to (5.1) $\mathbf{u}_T(t_{n-1})^e$ is defined on $\Gamma(t_n)$. The normal surface velocity \mathbf{w}_N is known, so a natural linearization of the nonlinear term in (5.2) is given by

$$\mathbf{P}\nabla\mathbf{u}_T^e(\mathbf{w}_N+\mathbf{u}_T)\approx\mathbf{P}(t_n)\nabla\mathbf{u}_T(t_n)^e(\mathbf{w}_N(t_n)+\mathbf{u}_T(t_{n-1})^e)\quad\text{on }\Gamma(t_n).$$

The FD approximations above need extensions of quantities defined on $\Gamma(t_j)$ to $\mathcal{O}(\Gamma(t_j))$. It is natural to consider a normal extension, which can be characterized as follows. Let $\mathbf{n} = \nabla \phi$ in $\mathcal{O}(\Gamma(t_j))$, where ϕ is the signed distance function for $\Gamma(t_j)$, and g a function defined on $\Gamma(t_j)$. The normal extension g^e of g is such that $g^e = g$ on $\Gamma(t_j)$ and

$$\mathbf{n} \cdot \nabla g^e = 0 \quad \text{in } \mathcal{O}(\Gamma(t_j)). \tag{5.3}$$

For practical purposes, ϕ can be a smooth level set function for Γ rather than a signed distance. In this case, the vector field $\mathbf{n} = \nabla \phi / |\nabla \phi|$ is normal on Γ but defines quasi-normal directions in a neighborhood. Extension of the velocity field along quasi-normal directions is equally admissible. We assume that at t = 0 an extension $\mathbf{u}_T(0)^e$ on $\mathcal{O}(\Gamma_0)$ solving (5.3) is given. We use the notation \mathbf{u}_T^j and p^j for an approximation of $\mathbf{u}_T(t_j)^e$ and $p(t_j)$, respectively. Based on the approximations above and (5.3) consider the following time discretization method for (2.12). Given $\mathbf{u}_T^0 = \mathbf{u}_T(0)^e$, for $n = 1, \ldots N$, find \mathbf{u}_T^n , defined on $\mathcal{O}(\Gamma(t_n))$ and tangential to $\Gamma(t_n)$, i.e. $(\mathbf{u}_T^n \cdot \mathbf{n})|_{\Gamma(t_n)} = 0$, and p^n defined on $\Gamma(t_n)$ such that

$$\begin{cases} \frac{\mathbf{u}_{T}^{n} - \mathbf{P}\mathbf{u}_{T}^{n-1}}{\Delta t} + \mathbf{P}\nabla\mathbf{u}_{T}^{n}(\mathbf{w}_{N}^{n} + \mathbf{u}_{T}^{n-1}) \\ + w_{N}^{n}\mathbf{H}\mathbf{u}_{T}^{n} - 2\mu\mathbf{P}\operatorname{div}_{\Gamma}E_{s}(\mathbf{u}_{T}^{n}) + \nabla_{\Gamma}p^{n} = \mathbf{f}^{n} \quad \text{on } \Gamma(t_{n}), \qquad (5.4) \\ \operatorname{div}_{\Gamma}\mathbf{u}_{T}^{n} = g^{n} \\ \mathbf{n} \cdot \nabla\mathbf{u}_{T}^{n} = 0 \quad \text{in } \mathcal{O}(\Gamma(t_{n})), \qquad (5.5) \end{cases}$$

with $w_N^n := w_N(t_n)$, $\mathbf{w}_N^n := \mathbf{w}_N(t_n)$, $\mathbf{f}^n := \mathbf{f}(t_n)$, $g^n := g(t_n)$. Geometric information in (5.4) is taken for $\Gamma(t_n)$, i.e. $\mathbf{P} = \mathbf{P}(t_n)$, $\mathbf{H} = \mathbf{H}(t_n)$. For space discretization, the stationary linearized surface PDE in (5.4) can be treated using a variational approach known from Refs. 16 and 18, in which the tangential constraint for the solution \mathbf{u}_T^n is relaxed using a penalty approach. This technique is now outlined. We set $\mathbf{c} := \mathbf{w}_N^n + \mathbf{u}_T^{n-1}$, $\Gamma_n := \Gamma(t_n)$ and introduce the following bilinear forms on Γ_n , with arguments \mathbf{u}, \mathbf{v} , vector functions on Γ_n that are not necessarily tangential:

$$A(\mathbf{u}, \mathbf{v}) = \frac{1}{\Delta t} \int_{\Gamma_n} \mathbf{u} \cdot \mathbf{P} \mathbf{v} \, ds + \int_{\Gamma_n} \mathbf{v} \cdot \mathbf{P}(\nabla \mathbf{u}) \mathbf{c} \, ds + \int_{\Gamma_n} w_N^n \mathbf{u}^T \mathbf{H} \mathbf{v} \, ds$$
$$+ 2\mu \int_{\Gamma_n} E_s(\mathbf{P} \mathbf{u}) : E_s(\mathbf{P} \mathbf{v}) \, ds + \tau \int_{\Gamma_n} u_N v_N \, ds, \qquad (5.6)$$

$$B(\mathbf{u}, p) = -\int_{\Gamma_n} p \operatorname{div}_{\Gamma} \mathbf{P} \mathbf{u} \, ds, \qquad (5.7)$$

where $\tau > 0$ is a penalty parameter. We introduce two Hilbert spaces

$$L_0^2(\Gamma_n) := \left\{ p \in L^2(\Gamma_n) \left| \int_{\Gamma_n} p \, ds = 0 \right\}, \text{ and} \right.$$
$$\mathbf{V}_* := \{ \mathbf{v} \in L^2(\Gamma_n)^3 \mid \mathbf{v}_T \in H^1(\Gamma_n)^3, v_N \in L^2(\Gamma_n) \},$$

with the norm $\|\mathbf{v}\|_{V_*}^2 = \|\mathbf{v}_T\|_{H^1(\Gamma_n)}^2 + \|v_N\|_{L^2(\Gamma_n)}^2$. A variational formulation corresponding to (5.4) is as follows: Find $\mathbf{u}_* \in \mathbf{V}_*$, $p \in L^2_0(\Gamma_n)$ such that

$$A(\mathbf{u}_*, \mathbf{v}) + B(\mathbf{v}, p) = \int_{\Gamma_n} \tilde{\mathbf{f}} \cdot \mathbf{v} \, ds \quad \text{for all } \mathbf{v} \in \mathbf{V}_*$$

$$B(\mathbf{u}_*, q) = -\int_{\Gamma_n} g^n q \, ds \quad \text{for all } q \in L^2_0(\Gamma_n),$$
(5.8)

with $\tilde{\mathbf{f}} := \mathbf{f} + \frac{1}{\Delta t} \mathbf{P} \mathbf{u}_T^{n-1}$. This variational formulation is consistent in the sense that if (\mathbf{u}_T^n, p^n) is a strong solution of (5.4) then $(\mathbf{u}_*, p) = (\mathbf{u}_T^n, p^n)$ solves (5.8). Using the Korn type inequality (3.11) and an inf-sup property of $B(\cdot, \cdot)$ it can be shown that for Δt sufficiently small and τ sufficiently large (but independent of Δt) the problem (5.8) is well-posed and its unique solution \mathbf{u}_* satisfies $\mathbf{u}_* \cdot \mathbf{n} = 0$, cf. Ref. 16 for a precise analysis. Therefore, for such Δt and τ Eq. (5.8) is a well-posed weak formulation of (5.4). For a finite element method introduced later it is important that the space \mathbf{V}_* admits vector functions not necessarily tangential to Γ_n . The solution \mathbf{u} of (5.8) is defined only on Γ_n and we do not specify an extension as in (5.5), yet. Such an extension will be determined in the finite element method, as explained in the following section.

Remark 5.1. In the practical implementation of a finite element method for (5.8) the surface $\Gamma(t_n)$ will be approximated by a piecewise planar surface Γ_h and the corresponding projection operator \mathbf{P}_h has discontinuities across boundaries between different planar segments of this approximate surface. This causes difficulties for the terms in the bilinear forms $A(\cdot, \cdot)$, $B(\cdot, \cdot)$ where derivatives of \mathbf{P}_h are involved. These can be avoided by eliminating derivatives of \mathbf{P} as follows. For $p \in H^1(\Gamma_n)$ we have $B(\mathbf{u}, p) = \int_{\Gamma_n} \nabla_{\Gamma} p \cdot \mathbf{Pu} \, ds = \int_{\Gamma_n} \nabla_{\Gamma} p \cdot \mathbf{u} \, ds$, which eliminates derivatives of \mathbf{P} . For the bilinear form $A(\cdot, \cdot)$ we can use the relation $\nabla_{\Gamma}(\mathbf{Pu}) = \nabla_{\Gamma}\mathbf{u} - u_N\mathbf{H}$ and thus $E_s(\mathbf{Pu}) = E_s(\mathbf{u}) - u_N\mathbf{H}$.

5.2. Finite element method

We now explain the spatial discretization of (5.8). Consider a fixed polygonal domain $\Omega \subset \mathbb{R}^3$ that strictly contains $\Gamma(t)$ for all $t \in (0,T)$. Let $\{\mathcal{T}_h\}_{h>0}$ be a family of shape-regular consistent triangulations of the bulk domain Ω , with $\max_{K \in \mathcal{T}_h} \operatorname{diam}(K) \leq h$. Corresponding to the bulk triangulation we define a standard finite element space of piecewise polynomial continuous functions of a fixed degree $k \geq 1$:

$$V_h^k = \{ v_h \in C(\Omega) : v_h \in P_k(K), \ \forall \ K \in \mathcal{T}_h \}.$$

$$(5.9)$$

The bulk velocity and pressure finite element spaces are standard Taylor–Hood spaces:

$$\mathbf{U}_h \coloneqq (V_h^2)^3, \quad Q_h \coloneqq V_h^1.$$

For efficiency reasons, we use an extension not in the given (*h* and Δt -independent) neighborhood $\mathcal{O}(\Gamma(t_n))$ of $\Gamma_n = \Gamma(t^n)$ but in a narrow band around Γ_n . This Δt -dependent narrow band consists of all tetrahedra within a δ_n distance from the surface, with

$$\delta_n := c_\delta \sup_{t \in (t_{n-1}, t_n)} \|w_N\|_{L^{\infty}(\Gamma(t))} \Delta t$$
(5.10)

and $c_{\delta} \geq 1$, an O(1) mesh-independent parameter. More precisely, we define the mesh-dependent narrow band

$$\mathcal{O}_{\Delta t}(\Gamma_n) := \bigcup \{ \overline{K} : K \in \mathcal{T}_h : \operatorname{dist}(\mathbf{x}, \Gamma_n) \le \delta_n \text{ for some } \mathbf{x} \in K \}.$$

We also need a subdomain of $\mathcal{O}_{\Delta t}(\Gamma_n)$ only consisting of tetrahedra intersected by Γ_n ,

$$\omega_{\Gamma}^{n} := \bigcup \{ \overline{K} \in \mathcal{T}_{h} : K \cap \Gamma_{n} \neq \emptyset \}.$$

In a time step from t_{n-1} to t_n , the surface may move up to $\Delta t \sup_{t \in (t_{n-1}, t_n)} \|w_N\|_{L^{\infty}(\Gamma(t))}$ distance in normal direction, which is thus the maximum distance from Γ_n to Γ_{n-1} . Therefore, c_{δ} in (5.10) can be taken sufficiently large, but independent of h, such that

$$\omega_{\Gamma}^n \subset \mathcal{O}_{\Delta t}(\Gamma_{n-1}). \tag{5.11}$$

This condition is the discrete analog of (5.1) and it is essential for the well-posedness of the finite element problem at time step n.

Next, we define finite element spaces for velocity and pressure as restrictions to the narrow band $\mathcal{O}_{\Delta t}(\Gamma_n)$ of the time-independent bulk spaces \mathbf{U}_h and Q_h :

$$\mathbf{U}_{h}^{n} := \{ \mathbf{v}|_{\mathcal{O}_{\Delta t}(\Gamma_{n})} \, | \, \mathbf{v} \in \mathbf{U}_{h} \}, \quad Q_{h}^{n} := \{ q|_{\mathcal{O}_{\Delta t}(\Gamma_{n})} \, | \, q \in Q_{h} \}.$$
(5.12)

Denote by $I_h(\mathbf{v}) \in \mathbf{U}_h^n$ the Lagrange interpolation of $\mathbf{v} \in C(\mathcal{O}_{\Delta t}(\Gamma_n))^3$. Our finite element formulation is based on formulation (5.8). Recall that in (5.8) we do not require \mathbf{u}_* to be tangential to $\Gamma(t)$. The tangential condition is weakly enforced by the penalty term in (5.6) with penalty parameter τ . Such a penalty approach is often used in finite element methods for vector values surface PDEs.^{14, 18, 27, 32} In the discretization in addition to this penalty term we include two volume terms with integrals over ω_{Γ}^n and $\mathcal{O}_{\Delta t}(\Gamma_n)$. The discrete problem is as follows: For given $\mathbf{u}_h^{n-1} \in \mathbf{U}_h^{n-1}$ and $\mathbf{c}_h^{n-1} = \mathbf{u}_h^{n-1} + I_h(w_N^e(t_n)\mathbf{n})$ find $\mathbf{u}_h^n \in \mathbf{U}_h^n$, $p_h^n \in Q_h^n$, satisfying

$$\int_{\Gamma_{n}} \left(\frac{\mathbf{u}_{h}^{n} - \mathbf{u}_{h}^{n-1}}{\Delta t} + (\nabla \mathbf{u}_{h}^{n}) \mathbf{c}_{h}^{n-1} + u_{N}^{n} \mathbf{H} \mathbf{u}_{h}^{n} \right) \cdot \mathbf{P} \mathbf{v}_{h} \, ds$$
$$+ 2\mu \int_{\Gamma_{n}} E_{s}(\mathbf{P} \mathbf{u}_{h}^{n}) : E_{s}(\mathbf{P} \mathbf{v}_{h}) \, ds + \tau \int_{\Gamma_{n}} (\mathbf{n} \cdot \mathbf{u}_{h}^{n}) (\mathbf{n} \cdot \mathbf{v}_{h}) \, ds$$
$$+ \int_{\Gamma_{n}} \nabla_{\Gamma} p^{n} \mathbf{v}_{h} \, ds + \rho_{u} \int_{\mathcal{O}_{\Delta t}(\Gamma_{n})} (\mathbf{n} \cdot \nabla \mathbf{u}_{h}^{n}) (\mathbf{n} \cdot \nabla \mathbf{v}_{h}) \, d\mathbf{x}$$
$$= \int_{\Gamma_{n}} \mathbf{f}^{n} \mathbf{v}_{h} \, ds \quad \forall \mathbf{v}_{h} \in \mathbf{U}_{h}^{n}$$
$$- \int_{\Gamma_{n}} \nabla_{\Gamma} q \, \mathbf{u}_{h}^{n} \, ds + \rho_{p} \int_{\omega_{\Gamma}^{n}} (\mathbf{n} \cdot \nabla p_{h}^{n}) (\mathbf{n} \cdot \nabla q_{h}) \, d\mathbf{x}$$
$$= \int_{\Gamma_{n}} g^{n} q_{h} \, ds \quad \forall q_{h} \in Q_{h}^{n}, \tag{5.13}$$

for n = 1, ..., N. The term $\int_{\mathcal{O}_{\Delta t}(\Gamma_n)} (\mathbf{n} \cdot \nabla \mathbf{u}_h^n) (\mathbf{n} \cdot \nabla \mathbf{v}_h) \, d\mathbf{x}$, with a parameter ρ_u , is included for two reasons. First, this term is often used in TraceFEM to improve the conditioning of the resulting stiffness matrix, e.g. see Ref. 7. Second, this volume term weakly enforces the extension condition (5.5) with $\mathcal{O}(\Gamma(t_n))$ replaced by $\mathcal{O}_{\Delta t}(\Gamma_n)$. In particular, at time t_n a well-conditioned algebraic system is solved for all discrete velocity degrees of freedom in the neighborhood $\mathcal{O}_{\Delta t}(\Gamma_n)$; we refer to Ref. 23 for a stability and convergence analysis of such an extension procedure for a scalar surface equation. The volume term in the pressure equation is added for the purpose of numerical stabilization of pressure.³¹ The formulation (5.13) is consistent in the sense that the equations hold if the solution of (5.4), extended along normal directions, is substituted instead of u_h^n and p_h^n . Penalty and stabilization parameters are set following the error analysis in Ref. 31:

$$\tau = h^{-2}, \quad \rho_u = h^{-1}, \quad \rho_p = h.$$

In practice, Γ_n , n = 1, 2, ..., is replaced by a sufficiently accurate approximation Γ_h^n in such a way that integrals over Γ_h^n can be computed accurately and efficiently. Other geometric quantities, i.e. n, H and P, are also replaced by sufficiently accurate approximations. The derivatives of projected fields, i.e. $E_s(\mathbf{Pu}_h^n)$ and $E_s(\mathbf{Pv}_h)$, are handled as discussed in Remark 5.1. For the surface Stokes problem discretized by the trace $\mathbf{P}_{k+1} - P_k$, k > 1, elements, the introduced geometric error is analyzed in Ref. 17. Below we will use the lowest order trace Taylor–Hood pair \mathbf{P}_2 – P_1 . An approximation Γ_h^n that is piece-wise planar with respect to \mathcal{T}_h leads to an $O(h^2)$ geometric error. This geometric error order is suboptimal given the interpolation order of the Taylor-Hood pair $\mathbf{P}_2 - P_1$. This suboptimality can be overcome by the isoparametric TraceFEM.¹² For numerical results in this paper, we use the following less efficient but simpler approach. For the geometry approximation (only) we construct a piece-wise planar Γ_{h}^{n} with respect to a local refinement of each tetrahedron from ω_{Γ}^n . The number of local refinement levels is chosen sufficiently large to restore the optimal $O(h^3)$ convergence. Note that this local refinement only influences the surface approximation and not the finite element spaces used.

Finally we note that the use of BDF2 instead of implicit Euler in the implicit time stepping scheme leads to obvious modifications of (5.13). In the experiments in the following section, we used this second order in time variant of (5.13).

6. Numerical Examples

For discretization, an initial triangulation \mathcal{T}_{h_0} was build by dividing $\Omega = (-\frac{5}{3}, \frac{5}{3})^3$ into 2^3 cubes and further splitting each cube into 6 tetrahedra with $h_0 = \frac{5}{3}$. Further, the mesh is refined in a sufficiently large neighborhood of a surface so that tetrahedra cut by $\Gamma(t)$ belong to the same refinement level for all $t \in [0, T]$. $\ell \in \mathbb{N}$ denotes the level of refinement and $h_{\ell} = \frac{5}{3} 2^{-\ell}$. The trace $\mathbf{P}_2 - P_1$ Taylor-Hood finite element method with BDF2 time stepping, as described in the previous section, is applied.

6.1. Convergence for a smooth solution

To verify the implementation and to check the convergence order of the discrete solution, we set up an experiment with a known tangential flow along an expanding/contracting sphere. In this example the total area of Γ is not preserved, but it allows to prescribe a flow **u** analytically and calculate **f** and *g*.

The surface Γ is given by its distance function

$$d(\mathbf{x},t) \coloneqq \|\mathbf{x}\| - r(t), \quad r(t) \coloneqq 1 + \frac{1}{4}\sin(2\pi t), \tag{6.1}$$

We consider $t \in [0, 1]$. The surface normal velocity is then $\mathbf{w}_N = w_N \mathbf{n}$, with $w_N(t) = r'(t) = \frac{\pi \cos(2\pi t)}{2}$, $\mathbf{n}(\mathbf{x}) = \mathbf{x}/|\mathbf{x}|$. We choose $\mu = 5 \times 10^{-3}$.

The exact solution is given by

$$\mathbf{u}(\mathbf{x},t) \coloneqq \mathbf{P}(\mathbf{x},t) \left(1 - 2t, 0, 0\right)^T, \quad p(\mathbf{x}) \coloneqq x y^2 + z, \tag{6.2}$$

and right-hand sides \mathbf{f} and $g = \operatorname{div}_{\Gamma} \mathbf{u}_T + w_N \kappa$ are computed accordingly from (6.1) and (6.2). For numerical integration, exact solutions and right-hand sides are extended along normal directions to Γ .

The numerical solution was computed on four consecutive meshes with refinement levels $\ell \in \{2, \ldots, 5\}$ and a time step $\Delta t = 0.05$ on level 2; Δt is halved in each spatial refinement, and for parameter in (5.10) we set $c_{\delta} = 1.1$. A mesh in $\mathcal{O}_{\Delta t}(\Gamma_n)$ together with the embedded $\Gamma(t)$ and computed solution is illustrated in Fig. 1. In Table 1, we show the mesh parameter h and the resulting (averaged over all time steps) number of active degrees of freedom (# d.o.f.). We see that a mesh refinement leads to approximately four times more degrees of freedom. Table 1 further reports the velocity and pressure errors measured in (approximate) $L_{V_1}^2$ and $L^2(S)$ norms.

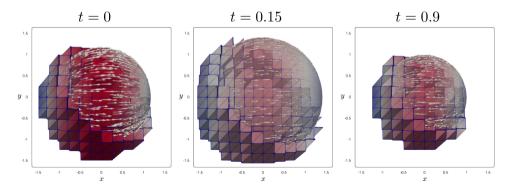


Fig. 1. Illustration of the extension mesh and solution at mesh level $\ell = 3$.

		2	$\frac{3}{2.08 \times 10^{-1}}$ 1.73×10^{4}		4	5
		17×10^{-1}			1.04×10^{-1}	5.21×10^{-2}
Averaged $\#$ d.	o.f. 4.	41×10^{3}	1.73 >	< 10*	6.82×10^{4}	2.73×10^5
$\ \mathbf{u} - \mathbf{u}_h\ _{1}$	Order	$\ \mathbf{u} - \mathbf{u}_h\ $	$ _{L^2(\mathcal{S})}$	Order	$ p - p_h _{L^2(S)}$	₅₎ Order
9.3×10^{-1}	1.3×1		0^{-1}		3.2×10^{-1}	
$1.9 imes 10^{-1}$	2.3	9.9×1	0^{-3}	3.72	$3.5 imes 10^{-2}$	3.2
$4.3 imes 10^{-2}$	2.13	9.2×1	0^{-4}	3.42	$7.3 imes 10^{-3}$	2.27
1.2×10^{-2}	1.92	1.2×1	0^{-4}	2.98	1.8×10^{-3}	2.02

Table 1. Convergence results for the example with analytical solution.

These norms were computed using a quadrature rule for time integration. Results demonstrate the expected second-order convergence in the "natural" norms and a higher order for the velocity error in the $L^2(S)$ norm. These orders are optimal for the \mathbf{P}_2-P_1 elements used.

6.2. Tangential flow on a deforming sphere

In this numerical example we consider a deforming unit sphere and compute the induced tangential flow, i.e. the numerical solution of the TSNSE (2.12). Denote by Γ_0 the reference sphere of radius 1 with the center in the origin O. Consider spherical coordinates $(r, \theta, \varphi) \in (0, \infty) \times [0, \pi] \times [0, 2\pi)$ and denote by $\mathcal{H}_n^m(\theta, \varphi)$, the spherical harmonic of degree n and order m. Assume that \mathcal{H}_n^m are normalized, i.e. $\|\mathcal{H}_n^m\|_{L^2(\Gamma_0)} = 1$. For the evolving surface we consider as ansatz

$$\Gamma(t) = \left\{ \mathbf{x} = (r, \theta, \varphi) \,|\, r = 1 + \sum_{n=1}^{N} \sum_{|m| \le n} A_{n,m}(t) \mathcal{H}_n^m(\theta, \varphi) \right\},\tag{6.3}$$

with suitably chosen coefficients $A_n(t)$. The function $\xi := \sum_{n=1}^N \sum_{|m| \le n} A_{n,m}(t)$ $\mathcal{H}_n^m(\theta, \varphi)$ describes the radial deformation. We assume small oscillations, $\|\xi\| \ll 1$. Under this assumption, an accurate *approximation* of the normal velocity is given by $\mathbf{w}_N = w_N \mathbf{n}$, with

$$w_N = \frac{d\xi}{dt} = \sum_{n=1}^N \sum_{|m| \le n} \frac{dA_{n,m}}{dt} \mathcal{H}_n^m, \quad \mathbf{n}(\mathbf{x}) = \mathbf{x}/|\mathbf{x}|.$$
(6.4)

We want the surface to be inextensible, i.e. $\frac{d}{dt}|\Gamma(t)| = 0$. Appropriate coefficients $A_{n,m}(t)$ such that we have inextensibility can be determined as follows. Application of the surface Reynolds transport formula and integration by parts gives for the variation of surface area:

$$\frac{d}{dt}|\Gamma(t)| = \frac{d}{dt}\int_{\Gamma(t)} 1\,ds = \int_{\Gamma(t)}\operatorname{div}_{\Gamma}\mathbf{w}_{N}\,ds = \int_{\Gamma(t)}\kappa w_{N}\,ds.$$
(6.5)

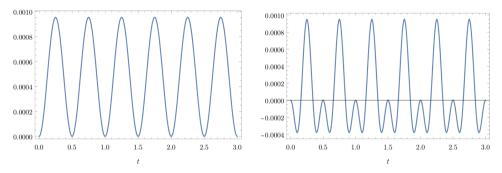


Fig. 2. Relative surface area variation $\frac{|\Gamma(0)| - |\Gamma(t)|}{|\Gamma(0)|}$ as a function of time for axisymmetric (left plot) and asymmetric (right plot) deformations of the sphere.

For the doubled mean curvature we have, cf. Ref. 21,

$$\kappa = 2 - 2\xi - \Delta_{\Gamma}\xi = 2 - \sum_{n=1}^{N} \sum_{|m| \le n} \{2A_{n,m}\mathcal{H}_{n}^{m} - n(n+1)A_{n,m}\mathcal{H}_{n}^{m}\}$$
$$= 2 + \sum_{n=1}^{N} \sum_{|m| \le n} A_{n,m}(n-1)(n+2)\mathcal{H}_{n}^{m}.$$

Using $\int_{\Gamma_0} \mathcal{H}_n^m = 0$, $n \ge 1$, and $\int_{\Gamma_0} \mathcal{H}_n^m \mathcal{H}_{n'}^{m'} = \delta_n^{n'} \delta_m^{m'}$, we compute for the area variation:

$$\frac{d}{dt}|\Gamma(t)| = \int_{\Gamma(t)} \kappa w_N ds = \sum_{n=1}^N \sum_{|m| \le n} (n-1)(n+2) \frac{dA_{n,m}}{dt} A_{n,m}.$$
 (6.6)

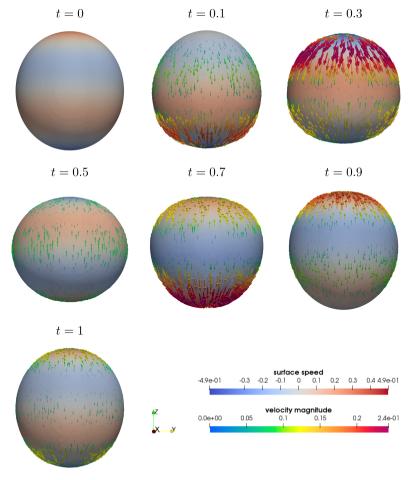


Fig. 3. Visualization of velocity field for axisymmetric deformations of the sphere; mesh level $\ell = 4$, $\Delta t = 0.01$. Click here to see the full animation.

Based on this formula we set $A_{2,0} = \frac{\epsilon}{2} \cos(\omega t)$, $A_{3,0} = \frac{\epsilon}{\sqrt{10}} \sin(\omega t)$, and $A_{n,m} = 0$ for other coefficients. For this choice of coefficients one easily verifies $\frac{d}{dt} |\Gamma(t)| = 0$. The TSNSE equations (2.12) are then solved with the right-hand side given by (2.13) with $u_N = w_N$ computed from (6.4). The initial velocity is zero.

In the first numerical example we let $\epsilon = 0.2$, $\omega = 2\pi$, $\mu = \frac{1}{2}10^{-4}$, and include \mathcal{H}_2^0 and \mathcal{H}_3^0 , two zonal spherical harmonics of degrees 2 and 3. The relative variation of the surface area $\Gamma(t)$ in the left plot in Fig. 2 shows less than 0.1% of surface variation. Non-zero variation is due to approximation errors and finite (rather than infinitesimal) deformations. The latter causes an approximation error in (6.4).

The velocity field induced by these axisymmetric deformations of the sphere is visualized in Fig. 3. We see that the velocity pattern is dominated by a

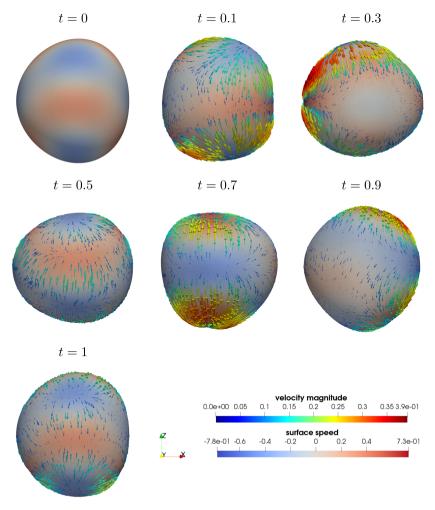


Fig. 4. Visualization of velocity field for asymmetric deformations of the sphere; mesh level $\ell = 4$, $\Delta t = 0.01$. Click here to see the full animation.

sink-and-source flow driven by the term $-\kappa w_N$ on the right-hand side of the divergence condition in (2.12).

We repeat the experiment, but decrease the viscosity to $\mu = \frac{1}{2}10^{-5}$ and add two more spherical harmonics, the sectorial \mathcal{H}_3^1 harmonic and the tesseral \mathcal{H}_4^2 one, to make the deformation asymmetric. The radial displacement in this experiment is then given by

$$\xi = 0.2 \left(\frac{1}{2} \cos(2\pi t) \mathcal{H}_2^0(\mathbf{x}) + \frac{1}{\sqrt{10}} \sin(2\pi t) \mathcal{H}_3^0(\mathbf{x}) \right), + 0.1 \left(\frac{1}{2} \cos(4\pi t) \mathcal{H}_3^1(\mathbf{x}) + \frac{5}{18} \sin(4\pi t) \mathcal{H}_4^2(\mathbf{x}) \right).$$

Again, the coefficients are such that $\frac{d}{dt}|\Gamma(t)| = 0$ according to Eq. (6.6). The resulting velocity field is visualized in Fig. 4. The velocity pattern is still dominated by the sink-and-source flow. Note that in both cases there are no outer forces and the flow is completely "geometry driven".

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References

- A. Alphonse, C. M. Elliott and B. Stinner, An abstract framework for parabolic PDEs on evolving spaces, *Portugal. Math.* **72** (2015) 1–46.
- V. Arnold, Sur la géométrie différentielle des groupes de lie de dimension infinie et ses applications à l'hydrodynamique des fluides parfaits, Ann. Inst. Fourier 16 (1966) 319–361.
- T. Aubin, Nonlinear Analysis on Manifolds. Monge-Ampere Equations (Springer, 1982).
- J. W. Barrett, H. Garcke and R. Nürnberg, Numerical computations of the dynamics of fluidic membranes and vesicles, *Phys. Rev. E* 92 (2015) 052704.
- A. Bonito, A. Demlow and M. Licht, A divergence-conforming finite element method for the surface Stokes equation, SIAM J. Numer. Anal. 58 (2020) 2764–2798.
- P. Brandner, A. Reusken and P. Schwering, On derivations of evolving surface Navier–Stokes equations, Interfaces and Free Boundaries (2022), https://doi.org/10.4171/ifb/482.
- E. Burman, P. Hansbo, M. G. Larson and A. Massing, Cut finite element methods for partial differential equations on embedded manifolds of arbitrary codimensions, *ESAIM: Math. Model. Numer. Anal.* 52 (2018) 2247–2282.
- C. Chan, M. Czubak and M. Disconzi, The formulation of the Navier–Stokes equations on Riemannian manifolds, J. Geom. Phys. 121 (2017) 335–346.
- G. Dziuk and C. Elliott, Finite elements on evolving surfaces, IMA J. Numer. Anal. 27 (2007) 262–292.

- G. Dziuk and C. M. Elliott, Finite element methods for surface PDEs, Acta Numerica 22 (2013) 289–396.
- D. G. Ebin and J. Marsden, Groups of diffeomorphisms and the motion of an incompressible fluid, Ann. Math. 92 (1970) 102–163.
- J. Grande, C. Lehrenfeld and A. Reusken, Analysis of a high-order trace finite element method for PDEs on level set surfaces, SIAM J. Numer. Anal. 56 (2018) 228–255.
- M. E. Gurtin and A. I. Murdoch, A continuum theory of elastic material surfaces, Arch. Ration. Mech. Anal. 57 (1975) 291–323.
- P. Hansbo, M. G. Larson and K. Larsson, Analysis of finite element methods for vector Laplacians on surfaces, *IMA J. Numer. Anal.* 40 (2020) 1652–1701.
- 15. P. Hartman, Ordinary Differential Equations, Vol. 590 (John Wiley and Sons, 1964).
- T. Jankuhn, M. A. Olshanskii and A. Reusken, Incompressible fluid problems on embedded surfaces: Modeling and variational formulations, *Interf. Free Bound.* 20 (2018) 353–377.
- T. Jankuhn, M. A. Olshanskii, A. Reusken and A. Zhiliakov, Error analysis of higher order trace finite element methods for the surface Stokes equations, *J. Numer. Math.* 29 (2021) 245–267.
- T. Jankuhn and A. Reusken, Higher order trace finite element methods for the surface Stokes equation, preprint (2019), arXiv:1909.08327.
- T. Jankuhn and A. Reusken, Trace finite element methods for surface vector-Laplace equations, *IMA J. Numer. Anal.* 41 (2020) 48–83.
- H. Koba, C. Liu and Y. Giga, Energetic variational approaches for incompressible fluid systems on an evolving surface, *Quart. Appl. Math.* 72 (2017) 359–389.
- 21. H. Lamb, Hydrodynamics (Univ. Press, 1924).
- 22. C. Lehrenfeld and M. Olshanskii, An Eulerian finite element method for PDEs in time-dependent domains, *ESAIM: M2AN* **53** (2019) 585–614.
- C. Lehrenfeld, M. A. Olshanskii and X. Xu, A stabilized trace finite element method for partial differential equations on evolving surfaces, *SIAM J. Numer. Anal.* 56 (2018) 1643–1672.
- T.-H. Miura, On singular limit equations for incompressible fluids in moving thin domains, Quart. Appl. Math. 76 (2018) 215–251.
- A. Murdoch and H. Cohen, Symmetry considerations for material surfaces, Arch. Ration. Mech. Anal. 72 (1979) 61–98.
- I. Nitschke, S. Reuther and A. Voigt, Hydrodynamic interactions in polar liquid crystals on evolving surfaces, *Phys. Rev. Fluids* 4 (2019) 044002.
- M. A. Olshanskii, A. Quaini, A. Reusken and V. Yushutin, A finite element method for the surface Stokes problem, *SIAM J. Sci. Comput.* 40 (2018) A2492–A2518.
- M. A. Olshanskii and A. Reusken, Trace finite element methods for PDEs on surfaces, in *Geometrically Unfitted Finite Element Methods and Applications*, eds. S. P. A. Bordas, E. Burman, M. G. Larson and M. A. Olshanskii (Springer International Publishing, 2017), pp. 211–258.
- M. Olshanskii, A. Reusken and J. Grande, A finite element method for elliptic equations on surfaces, SIAM J. Numer. Anal. 47 (2009) 3339–3358.
- M. A. Olshanskii, A. Reusken and X. Xu, An Eulerian space-time finite element method for diffusion problems on evolving surfaces, *SIAM J. Numer. Anal.* 52 (2014) 1354–1377.
- M. Olshanskii, A. Reusken and A. Zhiliakov, Inf-sup stability of the trace P2-P1 Taylor–Hood elements for surface PDEs, *Math. Comput.* 90 (2021) 1527–1555.
- M. A. Olshanskii and V. Yushutin, A penalty finite element method for a fluid system posed on embedded surface, J. Math. Fluid Mech. 21 (2019) 14.

- J. Prüss, G. Simonett and M. Wilke, On the Navier–Stokes equations on surfaces, J. Evol. Equ. 21 (2021) 3153–3179.
- P. Rangamani, A. Agrawal, K. K. Mandadapu, G. Oster and D. J. Steigmann, Interaction between surface shape and intra-surface viscous flow on lipid membranes, *Biomech. Model. Mechanobiol.* 12 (2013) 1–13.
- A. Reusken, Stream function formulation of surface Stokes equations, IMA J. Numer. Anal. 40 (2020) 109–139.
- S. Reuther and A. Voigt, The interplay of curvature and vortices in flow on curved surfaces, *Multiscale Model. Simul.* 13 (2015) 632–643.
- S. Reuther and A. Voigt, Erratum: The interplay of curvature and vortices in flow on curved surfaces, *Multiscale Model. Simul.* 16 (2018) 1448–1453.
- S. Reuther and A. Voigt, Solving the incompressible surface Navier–Stokes equation by surface finite elements, *Phys. Fluids* **30** (2018) 012107.
- P. Saffman and M. Delbrück, Brownian motion in biological membranes, Proc. Natl. Acad. Sci. 72 (1975) 3111–3113.
- G. Simonett and M. Wilke, H[∞]-calculus for the surface Stokes operator and applications, J. Math. Fluid. Mech. 24 (2022) 109.
- P. Steinmann, On boundary potential energies in deformational and configurational mechanics, J. Mech. Phys. Solids 56 (2008) 772–800.
- 42. M. E. Taylor, Analysis on Morrey spaces and applications to Navier–Stokes and other evolution equations, *Commun. Partial Differential Equations* **17** (1992) 1407–1456.
- R. Temam, Navier-Stokes Equations, Theory and Numerical Analysis (North-Holland, 1977).
- 44. R. Temam, Infinite-Dimensional Dynamical Systems in Mechanics and Physics (Springer, 1988).
- A. Torres-Sánchez, D. Millán and M. Arroyo, Modelling fluid deformable surfaces with an emphasis on biological interfaces, J. Fluid Mech. 872 (2019) 218–271.
- H. von Wahl, T. Richter and C. Lehrenfeld, An unfitted Eulerian finite element method for the time-dependent Stokes problem on moving domains, *IMA J. Numer. Anal.* 42 (2022) 2505–2544.
- A. Yavari, A. Ozakin and S. Sadik, Nonlinear elasticity in a deforming ambient space, J. Nonlinear Sci. 26 (2016) 1651–1692.