TRAVELING WAVE SOLUTIONS TO THE INCLINED OR PERIODIC FREE BOUNDARY INCOMPRESSIBLE NAVIER-STOKES EQUATIONS

JUNICHI KOGANEMARU AND IAN TICE

Abstract. This paper concerns the construction of traveling wave solutions to the free boundary incompressible Navier-Stokes system. We study a single layer of viscous fluid in a strip-like domain that is bounded below by a flat rigid surface and above by a moving surface. The fluid evolves in a domain that is bounded below by a flat rigid surface and above by a free moving surface that can be described by the graph of a continuous function, in dimensions \(n \geq 2\) (though, of course, only the cases \(n \in \{2, 3\}\) are physically relevant). Furthermore, we assume that the fluid evolves in a strip-like domain, bounded below by a flat, rigid, inclined surface and above by a free moving surface that can be described by the graph of a continuous function.

1. Introduction

The existence of traveling wave solutions to the equations of fluid mechanics has been a subject of intense study for nearly two centuries (see Section 1.4 for a brief summary). Until recently, most of the mathematical results in this area focused on inviscid fluids, but work in the last few years constructed traveling wave solutions to the free boundary Navier-Stokes equations with a single horizontally infinite but finite depth layer of incompressible fluid \([18]\), and with multiple layers \([25]\) in a uniform, downward-pointing gravitational field. The purpose of the present paper is to extend these constructions into more general physical configurations by considering two effects: inclination of the fluid domain, which results in a component of the gravitational field parallel to the fluid layer; and, periodicity of the fluid layer in certain directions. The key to the constructions in \([18, 25]\) was the identification and utilization of a new scale of anisotropic Sobolev spaces, which serve as the container space for the function describing the free surface of the fluid. The results in this paper rely crucially on some new functional analytic properties of these spaces, which we prove here: the development of this scale of spaces on domains with periodicity; and, the fact that these spaces form an algebra in the supercritical regime. While these results are essential for our specific PDE needs, they may be of independent interest to those interested in Sobolev spaces.

1.1. Kinematic and dynamic description of the fluid. In this paper we consider a single finite depth layer of viscous, incompressible fluid. We assume that the fluid evolves in a strip-like domain, bounded below by a flat, rigid, inclined surface and above by a free moving surface that can be described by the graph of a continuous function.

Prior to specifying the fluid domain or the equations of motion, we fix an orthogonal coordinate system as follows. We choose the unit vector \(e_n\) to be normal to the inclined surface, and we choose the unit vector \(e_1\) in such a way that \(\mathcal{G}\) lies in the \(e_1-e_n\) hyperplane. In other words, we posit that the uniform gravitational field \(\mathcal{G}\) resolves into two components as \(\mathcal{G} = \kappa e_1 - g e_n\), where \(\kappa \in \mathbb{R}\) and \(g \in (0, \infty)\). We can then define the angle of incline of the domain, \(\theta \in (-\pi/2, \pi/2)\), via \(\tan \theta = \kappa / g\) (see Figure 1).

Next we turn to a description of the fluid cross-section, \(\Sigma\), which will allow us to specify the free surface and rigid bottom of the fluid and to model periodicity. The flat rigid bottom of the fluid, which is orthogonal to the \(e_n\) direction, is described by \((n - 1)\)-dimensional sets of the form

\[
\Sigma = \begin{cases} \mathbb{R} \times \prod_{i=2}^{n-1} \Sigma_i, & \text{where } \Sigma_i = \mathbb{R} \text{ or else } \Sigma_i = L_i T, \\ \prod_{i=1}^{n-1} L_i T, & \end{cases}
\]

1.2. Mathematical preliminaries. A key component of our analysis is the development of some new functional analytic properties of a scale of anisotropic Sobolev spaces, including that these spaces are an algebra in the supercritical regime, which may be of independent interest.
where \( LT = \mathbb{R}/LZ \) denotes the 1-torus of periodicity length \( L > 0 \). This choice of \( \Sigma \) allows us to model: full periodicity in all directions, which corresponds to the second case; no periodicity, which corresponds to the first case with each \( \Sigma_i = \mathbb{R} \); or partial periodicity when \( n \geq 3 \), which corresponds to the first case with at least one \( \Sigma_i = L_iT \). However, for technical reasons that we will detail below, in the case of partial periodicity we cannot allow periodicity in the \( e_1 \) direction, and so the first factor of \( \Sigma \) must be \( \mathbb{R} \). To be more explicit in the physically relevant cases \( n \in \{2, 3\} \), we note that (1.1) allows for \( \Sigma \in \{\mathbb{R}, LT\} \) when \( n = 2 \), and when \( n = 3 \) it allows for \( \Sigma \in \{\mathbb{R}^2, \mathbb{R} \times L_T, L_1T \times L_2T\} \), but we exclude the possibility that \( \Sigma = L_1T \times \mathbb{R} \). In any dimension, we endow \( \Sigma \) with the usual topological and smooth structure.

With \( \Sigma \) in hand, we can now describe strip-like \( n \)-dimensional domains with cross section \( \Sigma \). Given any function \( \zeta : \Sigma \to (0, \infty) \), we define the set \( \Omega_\zeta \subset \Sigma \times \mathbb{R} \) via \( \Omega_\zeta = \{x = (x', x_n) \in \Sigma \times \mathbb{R} \mid 0 < x_n < \zeta(x')\} \) and the \( \zeta \)-graph surface \( \Sigma_\zeta \subset \Sigma \times \mathbb{R} \) via \( \Sigma_\zeta = \{x = (x', x_n) \in \Sigma \times \mathbb{R} \mid x_n = \zeta(x') \text{ for some } x' \in \Sigma\} \). We note that if \( \zeta \) is continuous, then the upper boundary of \( \Omega_\zeta \) is \( \Sigma_\zeta \) and its lower boundary is given by \( \Sigma_0 = \{x = (x', x_n) \in \Sigma \times \mathbb{R} \mid x_n = 0\} \).

We assume that at equilibrium, with all external forces and stresses absent, the fluid occupies the flat equilibrium domain \( \Omega_b = \{x = (x', x_n) \in \Sigma \times \mathbb{R} \mid 0 < x_n < b\} \) for some equilibrium depth parameter \( b > 0 \). Furthermore, we assume that under perturbation, the fluid occupies the time-dependent fluid domain \( \Omega_{b+\zeta}(t) \), where \( \zeta : \Sigma \times [0, \infty) \to (-b, \infty) \) is an unknown free surface function. The fluid is then bounded above by the \( (b + \zeta) \)-graph surface \( \Sigma_{b+\zeta}(t) \) and below by the flat boundary \( \Sigma_0 \).

In addition to the aforementioned gravitational force, we posit that there are four other distinct forces acting upon the fluid: one in the bulk and three on the free surface. The first bulk force is a generic force described by the vector field \( \mathbf{f}(\cdot, t) : \Omega_{b+\zeta}(t) \to \mathbb{R}^n \). The first surface force is a constant external pressure \( P_{ext} \in \mathbb{R} \) applied by the fluid above the free surface. The second surface force is generated by an externally applied stress tensor, which is described by a map \( \mathbf{T} : \Sigma_{b+\zeta}(t) \to \mathbb{R}_{sym}^{n \times n} \), where \( \mathbb{R}_{sym}^{n \times n} = \{M \in \mathbb{R}^{n \times n} \mid M = MT\} \) is the set of symmetric \( n \times n \) matrices. The symmetry condition is imposed to be consistent with the fact that stress tensors are typically symmetric in continuum mechanics, but this condition is not essential and could be dropped in our analysis. The third surface force is the surface tension generated by the surface itself, which we model in the standard way as \(-\sigma \mathcal{H}(\zeta)\), where \( \sigma > 0 \) is the coefficient of the surface tension, and

\[
\mathcal{H}(\zeta) = \text{div}'\left(\frac{\nabla \zeta}{\sqrt{1 + |\nabla \zeta|^2}}\right)
\]  

(1.2)

is the mean-curvature operator.

We assume that the evolution of the fluid is described for time \( t \geq 0 \) by its velocity field \( w(\cdot, t) : \Omega_{b+\zeta}(t) \to \mathbb{R}^n \) and its pressure \( P(\cdot, t) : \Omega_{b+\zeta}(t) \to \mathbb{R} \). For each \( t > 0 \), we require that the fluid velocity \( w \), the pressure \( P \), and

**Figure 1.** A sample portion of the time-dependent fluid domain under perturbation.
the free surface $\zeta$ satisfy the free boundary incompressible Navier Stokes equations:

$$
\begin{align*}
\rho (\partial_t w + w \cdot \nabla w) - \mu \Delta w + \nabla P &= -\rho g e_n + \rho \kappa e_1 + \hat{f}, & \text{in } \Omega_{b+\zeta(\cdot,t)} \\
\text{div } w &= 0, & \text{in } \Omega_{b+\zeta(\cdot,t)} \\
\partial_t \zeta &= w \cdot \nu \sqrt{1 + |\nabla \zeta|^2}, & \text{on } \Sigma_{b+\zeta(\cdot,t)} \\
(P I - \mu \nabla w) \nu &= [-\sigma H(\zeta) I + P_{ext} I + \bar{T}] \nu, & \text{on } \Sigma_{b+\zeta(\cdot,t)} \\
w &= 0, & \text{on } \Sigma_0.
\end{align*}
$$

(1.3)

Here $\rho > 0$ is the fluid density, $\mu > 0$ is the fluid viscosity, $\nabla w = (\nabla w) + (\nabla w)^T \in \mathbb{R}^{n \times n}$ is the symmetrized gradient of $w$, and $\nu(\cdot, t) = (-\nabla \zeta(\cdot, t), 1)/\sqrt{1 + |\nabla \zeta(\cdot, t)|^2} \in \mathbb{R}^n$ is the outward pointing unit normal to the surface $\Sigma_{b+\zeta(\cdot,t)}$. The first two equations in (1.3) are the standard incompressible Navier-Stokes equations; the first equation asserts the Newtonian balance of forces, the second enforces the conservation of mass. The third equation is the kinematic boundary condition describing the evolution of the free surface with the fluid. We note in particular that the third equation may be written as a transport equation in the form of $\partial_t \zeta + \nabla' \zeta \cdot w'|\Sigma_{b+\zeta(\cdot,t)} = w_n|\Sigma_{b+\zeta(\cdot,t)}$, which shows that the free surface $\zeta$ is transported by the horizontal component of the velocity $w'$ and driven by the vertical component of the velocity $w_n$. The fourth equation encodes the dynamic boundary conditions asserting the balance of forces on the free surface. The fifth equation is the typical no-slip condition enforced on flat rigid surfaces.

For the sake of convenience, we shall assume without loss of generality that $\rho = \mu = g = 1$. This can be achieved by dividing both sides of the first equation in (1.3) by $\rho$, rescaling the space-time variables, and renaming $b, \sigma, \kappa, \hat{f}, P_{ext}, \bar{T}$, and (in the periodic settings) the periodicity length scales $L_i$, accordingly. This yields the system

$$
\begin{align*}
\partial_t w + w \cdot \nabla w - \Delta w + \nabla P &= -e_n + \kappa e_1 + \hat{f}, & \text{in } \Omega_{b+\zeta(\cdot,t)} \\
\text{div } w &= 0, & \text{in } \Omega_{b+\zeta(\cdot,t)} \\
\partial_t \zeta &= w \cdot \nu \sqrt{1 + |\nabla \zeta|^2}, & \text{on } \Sigma_{b+\zeta(\cdot,t)} \\
(P I - \nabla w) \nu &= [-\sigma H(\zeta) I + P_{ext} I + \bar{T}] \nu, & \text{on } \Sigma_{b+\zeta(\cdot,t)} \\
w &= 0, & \text{on } \Sigma_0.
\end{align*}
$$

(1.4)

For a differentiable vector field $u$ and a scalar $p$, we define the stress tensor $S(p, u) = pI - \nabla u \in \mathbb{R}^{n \times n}$, where $I$ is the $n \times n$ identity and $\nabla u$ is the symmetrized gradient of $u$. By defining the divergence operator to act on tensors in the canonical way, we have $\text{div } S(p, u) = \nabla P - \Delta u + \nabla \text{div } u$. This means that in the first equation of (1.4) we can rewrite

$$
\nabla P - \Delta u = \text{div } S(P, w).
$$

(1.5)

1.2. Shear flows and perturbations. The system (1.4) admits steady shear flow solutions that reduce to hydrostatic equilibrium when $\kappa = 0$. To see this, we suppose that $\zeta = 0, \hat{f} = 0, \bar{T} = 0$. We then define the smooth functions $s_0, s : \mathbb{R} \rightarrow \mathbb{R}$ via

$$
s_0(x_n) = bx_n - \frac{x_n^2}{2} \quad \text{and} \quad s(x_n) = \kappa s_0(x_n).
$$

(1.6)

We then define the steady shear velocity field $U^1 : \Omega_b \rightarrow \mathbb{R}^n$ via

$$
W^1(x', x_n) = s(x_n)e_1,
$$

(1.7)

and the equilibrium hydrostatic pressure $p : \Omega_b \rightarrow \mathbb{R}$ via

$$
p(x', x_n) = P_{ext} - x_n + b.
$$

(1.8)

See Figure 1 for a sketch of the $W^1$ profile. One can readily check that $w = W^1, P = p, \zeta = 0$ is a steady shear flow solution to (1.4) when $\hat{f} = 0, \bar{T} = 0$. Note that these shear flow solutions are special solutions induced by $\kappa \neq 0$, and they exist due to the presence of viscosity, with no clear analogue in the Euler system. In the literature, these solutions are sometimes referred to as Nusselt solutions.

We will study the system (1.4) as a perturbation of this steady solution. We define the perturbation of the velocity field and pressure field given by

$$
\bar{w}^1(x, t) = w(x, t) + W^1(x), \quad \bar{P}^1(x, t) = P(x, t) + p(x),
$$

(1.9)
and we see that $(\bar{w}^1, \bar{T}^1, \zeta)$ is a solution to (1.4) if and only if $(w, P, \zeta)$ satisfies

$$
\begin{align*}
\partial_t w + \text{div} S(P, w) + (w + W^1) \cdot \nabla (w + W^1) &= \bar{f}, \\
\text{div} w &= 0, \\
\partial_t \zeta + \nabla \zeta \cdot w' + (\partial_t \zeta) s(\zeta + b) &= w_n, \\
S(P, w) \nu &= [(\zeta - \sigma \mathcal{H}(\zeta)) I + \bar{T} - \kappa \zeta (e_1 \otimes e_n + e_n \otimes e_1)] \nu, \\
w &= 0,
\end{align*}
$$

in $\Omega_{b+\zeta(\cdot, \cdot)}$, in $\Omega_{b+\zeta(\cdot, \cdot)}$, on $\Sigma_{b+\zeta(\cdot, \cdot)}$, on $\Sigma_{b+\zeta(\cdot, \cdot)}$, on $\Sigma_0$, respectively. (1.10)

Here we have utilized the identity (1.5) in the first equation of (1.10), and the tensor product $v \otimes w \in \mathbb{R}^n \times \mathbb{R}^n$ of two vectors $v, w \in \mathbb{R}^n$ is defined in the standard way via $(v \otimes w)_{ij} = v_i w_j$.

### 1.3. Traveling wave solutions around shear flows

In this paper our main goal is to construct traveling wave solutions to (1.10), which are solutions that are stationary when viewed in a coordinate system moving at a constant speed. For this stationary condition to hold, the moving coordinate system must travel at a constant speed parallel to the flat rigid surface $\Sigma$. In this paper we assume that the traveling waves move at a constant speed in the direction of incline; in other words, that the moving coordinate system’s velocity relative to the Eulerian coordinates of (1.3) is $\gamma e_1$ for some $\gamma \in \mathbb{R} \setminus \{0\}$. The speed of the traveling wave is then $|\gamma|$, and $\text{sgn}(\gamma)$ indicates the direction of travel along the $e_1$ axis.

Next we proceed to reformulate the problem in the traveling coordinates. We assume that the stationary free surface is given by an unknown function $\eta : \Sigma \to (-b, \infty)$, and it is related to $\zeta$ via $\zeta(x', t) = \eta(x' - \gamma t e_1)$. The stationary velocity, pressure, force, and stress $v, q, \nabla, \mathcal{T}$ are related to $w, P, \bar{f}, \bar{T}$ via

$$
w(x, t) = v(x - \gamma t e_1), \quad P(x, t) = q(x - \gamma t e_1), \quad \bar{f}(x, t) = f(x - \gamma t e_1), \quad \bar{T}(x, t) = T(x' - \gamma t e_1).$$

(1.11)

Then using (1.11), the system (1.10) reduces to a time-independent system for $(v, q, \eta)$ given the forcing terms $f$ and $\mathcal{T}$,

$$
\begin{align*}
\text{div} S(q, v) - \gamma e_1 \cdot \nabla v + (v + W^1) \cdot \nabla (v + W^1) &= f, \\
\text{div} v &= 0, \\
-\gamma \partial_t \eta + \nabla \eta \cdot v' + (\partial_t \eta) s(\eta + b) &= v_n, \\
S(q, v) \mathcal{N} &= [(\eta - \sigma \mathcal{H}(\eta)) I + \mathcal{T} - \kappa \eta (e_1 \otimes e_n + e_n \otimes e_1)] \mathcal{N}, \\
v &= 0,
\end{align*}
$$

in $\Omega_{b+\eta}$, in $\Omega_{b+\eta}$, on $\Sigma_{b+\eta}$, on $\Sigma_{b+\eta}$, on $\Sigma_0$, respectively. (1.12)

where we have defined the non-unit normal vector field

$$
\mathcal{N} = (-\nabla' \eta, 1).
$$

(1.13)

For technical reasons to be discussed later in Section 1.7, it is convenient for us to remove the $\eta$ terms appearing in the fourth equation of (1.12) by introducing an additional time-independent perturbation term

$$
W^2(x', x_n) = \kappa x_n \eta(x') e_1
$$

(1.14)

and by defining the modified shear velocity and modified pressure via

$$
\bar{v}^1(x, t) = v + W^2(x), \quad \bar{q}^1(x) = q(x) + \eta(x').
$$

(1.15)

One may readily check that $(\bar{v}^1, \bar{q}^1, \zeta)$ is a solution to (1.12) if and only if $(v, q, \eta)$ satisfies

$$
\begin{align*}
\text{div} S(q, v) - \gamma e_1 \cdot \nabla v - \kappa x_n \partial_n \eta e_1 + (v + W^1 + W^2) \cdot \nabla (v + W^1 + W^2) + (\nabla' \eta, 0) &= f, \\
-\kappa x_n \nabla' \eta (x') e_1 - \kappa (x_n \nabla \partial_1 \eta, \partial_1 \eta) &= 0, \\
\text{div} v + \kappa x_n \partial_1 \eta e_1 &= 0, \\
-\gamma + s(\eta + b) + \kappa (\eta + b) \eta \partial_1 \eta &= v \cdot \mathcal{N}, \\
S(q, v) \mathcal{N} &= [-\sigma \mathcal{H}(\eta) I + \mathcal{T} + \kappa (b + \eta) (e_1 \otimes (\nabla' \eta, 0) + (\nabla' \eta, 0) \otimes e_1)] \mathcal{N}, \\
v &= 0,
\end{align*}
$$

in $\Omega_{b+\eta}$, in $\Omega_{b+\eta}$, on $\Sigma_{b+\eta}$, on $\Sigma_{b+\eta}$, on $\Sigma_0$, respectively. (1.16)

We note that by using the perturbations introduced in (1.15), we have replaced the $\eta$ terms in the fourth equation of (1.12) by either derivatives of $\eta$ or products of $\eta$ and its derivatives, at the price of introducing additional terms in the first equation.
1.4. Previous work. The system (1.3) and its variants have been studied extensively in the mathematical literature. For a brief survey of the subject, we refer to Section 1.2 of Leoni-Tice [18] and Section 1.4 of Tice [28]. For a more thorough review of the subject, we refer the surveys of Tolland [29], Groves [10], Strauss [26], the recent paper by Strauss et al. [15] in the inviscid case, and to the surveys of Zadrzyńska [32] and Shibata-Shimizu [23] in the viscous case.

When $\kappa = 0$, the small data theory for the free boundary problem (1.3) over periodic domains is well-understood in dimension $n = 3$. For the problem with surface tension, Nishida-Teramoto-Yoshihara [20] constructed global periodic solutions and proved that they decay exponentially fast to equilibrium. For the problem without surface tension, Hataya [14] constructed global solutions decaying at a fixed algebraic rate, and later Guo-Tice [12] constructed global solutions decaying almost exponentially. In the non-periodic setting, Beale [3] established the local-wellposedness of solutions without surface tension, and Beale [4] proved the global existence of solutions with surface tension. Beale-Nishida [5] later established that the aforementioned solutions decay at an algebraic rate. Tani-Tanaka [27] proved the global existence solutions with and without surface tension under milder assumptions on the initial data, but did not study their decay rates. Guo-Tice [13] proved that for the problem without surface tension, the global solutions decay at a fixed algebraic rate.

The investigation of the dynamics of viscous shear flows without free boundary is classical and dates back to the work of Orr [21] and Sommerfeld [24], where they noted the so-called viscous destabilization phenomenon. This was subsequently investigated formally by many authors in the physics literature, including Heisenberg [16], Lin [19], and Tollmien [30]. However, it wasn’t until recently that a rigorous mathematical proof for the instability of viscous incompressible shear flows appeared, in the work of Grenier-Guo-Nguyen [9].

When $\kappa \neq 0$, much less is known about the dynamics of the free boundary problem (1.3). Ueno [31] studied the 2D problem with surface tension in the thin film regime and established uniform estimates of solutions with respect to the thinness parameter. Padula [22] studied the 3D problem with surface tension and proved sufficient conditions for asymptotic stability under a priori assumptions on the global existence of solutions. Tice [28] studied the asymptotic stability of shear flow solutions to the nonlinear problem with and without surface tension, and proved that solutions decay exponentially fast to equilibrium with surface tension and almost exponentially without surface tension.

The existence of traveling wave solutions without background shear flows to (1.3) first appeared in the recent work of Leoni-Tice [18], and the results therein were extended by Stevenson-Tice [25] to a multilayer configuration. To the best of our knowledge, there are no known results on the existence of traveling wave solutions around shear flows in the case when $\kappa \neq 0$ or when the cross section $\Sigma$ is periodic.

1.5. Reformulation in a fixed domain. We note that in (1.16), the fluid domain of interest $\Omega_{b+\eta}$ is dependent on an unknown free surface $\eta$. To bypass the difficulty of working in such a domain, we proceed to reformulate the problem on a fixed domain $\Omega_b = \Sigma \times (0, b)$, at the cost of worsening the nonlinearities in the system.

To do so we introduce the flattening map $\mathfrak{F}_\eta : \Omega_b \to \overline{\Omega_{b+\eta}}$ associated to a continuous function $\eta : \Sigma_b \to (-b, \infty)$, defined via

$$\mathfrak{F}_\eta(x', x_n) = x + \frac{x_n \eta(x')}{b} e_n.$$  

(1.17)

We note that by construction, $\mathfrak{F}_\eta|\Sigma_0 = \Sigma_0$ and $\mathfrak{F}_\eta(x', b) = \Sigma_{\eta+b}$ (see Figure 2). Moreover, $\mathfrak{F}_\eta$ is bijective with its inverse $\mathfrak{F}_\eta^{-1} : \Omega_{b+\eta} \to \Omega_b$ given by

$$\mathfrak{F}_\eta^{-1}(y) = \left(y', \frac{b y_n}{b + \eta(x')} \right).$$  

(1.18)

![Figure 2. The flattening map $\mathfrak{F}_\eta$ and its inverse](image-url)
Throughout the rest of the paper we will typically suppress the dependence of the subsequent maps and their associated domains on $b$ and $\eta$, e.g. writing $\mathfrak{F}_{\eta}$ as $\mathfrak{F}$ and $\Omega_{b}$ as $\Omega$, unless these dependencies need to be emphasized. We also use the following abuse of notation throughout the paper: since the hypersurface $\Sigma_{b} \subset \Sigma \times \mathbb{R}$ is canonically diffeomorphic to $\Sigma$ via the projection map $\pi: \Sigma_{b} \to \Sigma$ given by $\pi(x',b) = x'$, we will use this to identify $H^{s}(\Sigma_{b}; V)$ with $H^{s}(\Sigma; V)$ (and similar spaces) for any finite dimensional vector space $V$. This allows us to sometimes write $f(x')$ instead of $f(x', b)$, and allows us to apply the horizontal Fourier transform on $\Sigma_{b}$ in the natural way.

Following this convention, we note that $\mathfrak{F}$ is a homeomorphism inheriting the regularity of $\eta$, in the sense that if $\eta \in C^{k}(\Sigma; \mathbb{R})$ then $\mathfrak{F}$ is a $C^{k}$ diffeomorphism. When $\eta$ is differentiable we may compute

$$\nabla \mathfrak{F}(x) = \begin{pmatrix} I_{(n-1)\times(n-1)}(x, \eta(x')) & 0_{(n-1)\times1} \\ \frac{-x_{n}\eta'(x')} {b + \eta(x')} & 1 \end{pmatrix} + \begin{pmatrix} 0_{(n-1)\times1} \\ \frac{-x_{n}\eta'(x')} {b + \eta(x')} \end{pmatrix},$$

(1.19)

Thus, the Jacobian $J$ and the inverse Jacobian $K$ of $\mathfrak{F}$ are

$$J = \det \nabla \mathfrak{F} = 1 + \frac{\eta} {b + \eta}, \quad K = \frac{1} {J} = \frac{b} {b + \eta}. \quad (1.20)$$

We then introduce the matrix $A: \Omega_{b} \to \mathbb{R}^{n\times n}$ defined via

$$A(x) = (\nabla \mathfrak{F})^{-\top} = \begin{pmatrix} I_{(n-1)\times(n-1)}(x, \eta(x')) & 0_{(n-1)\times1} \\ \frac{-x_{n}\eta'(x')} {b + \eta(x')} & 1 \end{pmatrix} = \begin{pmatrix} I_{(n-1)\times(n-1)} & 0_{(n-1)\times1} \\ \frac{-x_{n}\eta'(x')} {b + \eta(x')} & 1 \end{pmatrix},$$

(1.21)

the $A$-dependent differential operators

$$(\nabla_{A}f)_{ij} = \sum_{j=1}^{n} A_{ij}\partial_{j}f, \quad (X \cdot \nabla_{A}u)_{i} = \sum_{j,k=1}^{n} X_{j}A_{jk}\partial_{k}u_{i}, \quad \text{div}_{A}X = \sum_{i,j=1}^{n} A_{ij}\partial_{j}X_{i},$$

and also

$$(\mathbb{D}_{A}u)_{ij} = \sum_{k=1}^{n} A_{ik}\partial_{k}u_{j} + A_{jk}\partial_{k}u_{i}, \quad S_{A}(p,u) = pI - \mathbb{D}_{A}u,$$

$$\text{div}_{A}S_{A}(p,u) = \nabla_{A}p - \Delta_{A}u - \nabla_{A} \text{div}_{A}u, \quad (\Delta_{A}u)_{i} = \sum_{j,k=1}^{n} A_{ij,k}\partial_{k}(A_{jm}\partial_{m}u_{i}). \quad (1.23)$$

Now writing $u, U^{1}, U^{2}: \Omega \to \mathbb{R}^{n}$, $p: \Omega \to \mathbb{R}$, $f: \Omega \to \mathbb{R}^n$ via $u = v \circ \mathfrak{F}, U^{1} = W^{1} \circ \mathfrak{F}, U^{2} = W^{2} \circ \mathfrak{F}, p = q \circ \mathfrak{F}$, where $v, q$ satisfies (1.10), (1.16) can be reformulated as the quasilinear system

$$\begin{cases}
\text{div}_{A}S_{A}(p,u) - \gamma_{1} \cdot \nabla_{A}u - \gamma_{2}(x_{n} + \eta \frac{x_{n}}{b}) \partial_{n} \eta_{1} + (u + U^{1} + U^{2}) \cdot \nabla_{A}(u + U^{1} + U^{2}) \\
+ (\nabla_{A}u, 0) - \kappa(x_{n} + \eta \frac{x_{n}}{b}) \Delta\eta_{1} - \kappa((x_{n} + \eta \frac{x_{n}}{b}) \nabla\partial_{1}\eta, \partial_{1}\eta) = f \circ \mathfrak{F}_{\eta}, & \text{in } \Omega \\
\text{div}_{A}u + \kappa(x_{n} + \eta \frac{x_{n}}{b}) \partial_{n} \eta = 0, & \text{in } \Omega \\
u \cdot \mathcal{N} = (-\gamma + s(\eta + b) + \kappa(\eta + b) \eta) \partial_{1}\eta, & \text{on } \Sigma_{b} \\
S_{A}(p,u)\mathcal{N} = [-\sigma\mathcal{H}(\eta)I + T \circ \mathfrak{F}_{\eta} + \eta \partial_{1}\eta + \kappa(\eta + b)(e_{1} \otimes (\nabla\eta, 0) + (\nabla\eta, 0) \otimes e_{1})] \mathcal{N}, & \text{on } \Sigma_{b} \\
u = 0, & \text{on } \Sigma_{0}.
\end{cases}\quad (1.24)$$

1.6. Statement of main results. We now state the main results obtained in this paper, though we delay a thorough discussion of them to the next subsection. The first result concerns a key property of the anisotropic Sobolev space $X^{s}(\Sigma; \mathbb{F})$, defined later in Section 2.1 for $\mathbb{F} \in \{\mathbb{R}, \mathbb{C}\}$, which serves as the container space for the free surface function $\eta$ when $\mathbb{F} = \mathbb{R}$.

**Theorem 1.1** (Proven later in Section 2.2). Suppose $s > d/2$, where $d = \dim \Sigma$ and $\Sigma$ is defined via (1.1). Then the anisotropic space $X^{s}(\Sigma; \mathbb{F})$ defined in Definition 2.5 is an algebra.

The proof of this theorem utilizes anisotropic Littlewood-Paley techniques heavily inspired by [11]. We note that this theorem is only non-trivial in the case when $1 \in R_{2}; |R_{2}| \geq 1$, where the set $R_{2}$ is defined via (A.2). It turns out that in all the remaining possibilities of $\Sigma$, the anisotropic Sobolev $X^{s}(\Sigma; \mathbb{F})$ coincides with the standard Sobolev space $H^{s}(\Sigma; \mathbb{F})$, so the result follows directly from standard Sobolev theory. In Lemma 2.2 we give the precise characterization of product domains $\Gamma$ for which $H^{s}(\Gamma; \mathbb{F}) = X^{s}(\Gamma; \mathbb{F})$ and $H^{s}(\Gamma; \mathbb{F}) \subseteq X^{s}(\Gamma; \mathbb{F})$. Moreover, we show in Proposition 2.3 that if $\Gamma = L_{1}\mathbb{T} \times \mathbb{R}$, then the space $X^{s}(\Gamma; \mathbb{F})$ is not complete, which is an initial indication of why we cannot allow for $\Sigma = L_{1}\mathbb{T} \times \mathbb{R}$ in our analysis.

With Theorem 1.1 in hand, we are able to prove the solvability of flattened system (1.24), and by extension the solvability of the original system (1.16). Before stating the solvability results we make a quick comment on the
forcing terms in (1.16) and (1.24). We first note that the bulk force \( f \) in (1.16) needs to be defined in the domain \( \Omega_{b+n} \), which depends on the unknown free surface function \( \eta \). In order to work with \( f \) without the implicit need to know \( \eta \) first, we will assume a priori that \( f \) is defined in all of \( \Sigma \times \mathbb{R} \). We will also need the map \((f, \eta) \mapsto \mathfrak{F}_\eta\) to be \( C^1 \), since we need to invoke the implicit function theorem in our analysis in Section 4. It is known that this is possible, for example, if the domain of \( f \) is \( H^{s+1} \) and its codomain is \( H^s \). Due to the strip-like structure of the fluid domain and the flattening map defined in (1.17), we can in fact allow for another type of bulk force which preserves the regularity of its domain. A detailed discussion of this can be found in Section 1.4 of [18]. Moving forward, we consider the generalized bulk forces \( f + L_{\Omega_{b+n}} f \) and the generalized surface forces \( T|_{\Sigma_{b+n}} + s_{b+n} T \) in (1.16), where \( L_{\Omega_{b+n}} f(x', x_n) = f(x') \) and \( s_{b+n} T(x', x_n) = T(x') \).

With this in mind, we are able to prove the solvability of (1.24). The spaces \( C_k^b, C_0^k, H^s(\Omega; \mathbb{R}^n) \) referred to in the following results are defined in Section 1.8, and the space \( X^s \) is defined in Definition 3.10. We also write \( \Sigma \ni x' = x'_{R\Sigma} + x'_{T\Sigma} \) in accordance with the notation introduced in (A.1).

**Theorem 1.2** (Proved later in Section 4.2). Let \( \Sigma \) be given by (1.1) and suppose that \( n \equiv s > \frac{n}{2} \). Further suppose that either \( \sigma > 0 \) and \( n \geq 2 \) or else \( \sigma = 0 \) and \( n = 2 \). Then there exists open sets

\[
U^s \subset (\mathbb{R} \setminus \{0\}) \times \mathbb{R} \times H^{s+2}(\Sigma \times \mathbb{R}; \mathbb{R}^{n \times n}) \times H^{s+\frac{3}{2}}(\Sigma; \mathbb{R}^{n \times n}) \times H^{s+1}(\Sigma \times \mathbb{R}; \mathbb{R}^n) \times H^s(\Sigma; \mathbb{R}^n)
\]

and \( O^s \subset X^s \) such that the following hold.

1. \((0,0,0) \in O^s\), and for every \((u,p,\eta) \in O^s\) we have that

\[
u \in C^2_b^{s-\frac{n-2}{2}}(\Omega; \mathbb{R}^n), p \in C^1_b^{s-\frac{n-2}{2}}(\Omega; \mathbb{R}), \eta \in C^0_b^{s-\frac{n-2}{2}}(\Sigma; \mathbb{R})
\]

and

\[
limit_{|x'_{R\Sigma} \to \infty} \partial^\alpha u(x) = \theta \text{ for all } \alpha \in \mathbb{N}^n \text{ such that } |\alpha| \leq 2 + \left[ s - \frac{n}{2} \right]
\]

2. We have \((\mathbb{R}^n \setminus \{0\}) \times \{0\} \times \{0\} \times \{0\} \times \{0\} \subset U^s\)

3. For each \((\gamma, \kappa, \tau, T, f, f) \in U^s\), there exists a unique \((u, p, \eta) \in O^s\) classically solving

\[
\begin{align*}
\text{div}_\mathcal{A} S_\mathcal{A}(p, u) &= -\gamma e_1 \cdot \nabla \eta \cdot \partial_1 \eta + e_1 + (u + U^1 + U^2) \cdot \nabla \eta + \kappa((x_n + \frac{\eta}{\tau^2}) \nabla^\gamma \partial_1 \eta, \partial_1 \eta) = f \circ \mathfrak{F}_\eta + L_{\Omega_{b+n}}, &\text{in } \Omega \\
\text{div}_\mathcal{A} u + \kappa((x_n + \frac{\eta}{\tau^2}) \nabla^\gamma \partial_1 \eta) &= 0, &\text{in } \Omega \\
u \cdot \nabla \nabla^\gamma \partial_1 \eta &= 0, &\text{on } \Sigma_b \\
S_\mathcal{A}(p, u) \nabla^\gamma \partial_1 \eta &= -\sigma \mathcal{H}(\eta) I + T \circ \mathfrak{F}_{\Sigma_b} + s_{b+n} T + s_{\kappa}((e_1 \nabla^\gamma \eta, 0) + (\nabla^\gamma \eta, 0) \otimes e_1), &\text{on } \Sigma_b \\
\gamma &= 0, &\text{on } \Sigma_0.
\end{align*}
\]

4. The map \( U^s \ni (\gamma, \kappa, \tau, T, f, f) \mapsto (u, p, \eta) \in O^s \) is \( C^1 \) and locally Lipschitz.

A few remarks are in order. First, we note again that in the physically relevant case \( n = 3 \), we cannot solve (1.29) in the stated spaces if \( \Sigma = L_T \times \mathbb{R} \). Second, Theorem 1.2 is analogous to Theorem 1.1 in [18] for the problem without incline and periodicity, and Theorem 1 in [25] for the multilayer problem. However, our choice of the function space \( X^s \) is slightly different from the one used there because we have formulated the problem (1.29) in a slightly different manner, with the \( \eta \) term shifted from the dynamic boundary condition into the bulk. This results in the pressure belonging to a standard Sobolev space rather than a specialized anisotropic one, as in [18, 25]. Third, we can say something about the set of \((\gamma, \kappa)\) parameters for which we can produce solutions,

\[
\mathfrak{P}^s = \{(\gamma, \kappa) \in \mathbb{R} \setminus \{0\} \times \mathbb{R} \mid (\gamma, \kappa, \tau, T, f, f) \in U^s \text{ for some } (T, f, f) \}.
\]

Indeed, an examination of our proof of Theorem 1.2 shows that for every \( \gamma \in \mathbb{R} \setminus \{0\} \), there exists a \( \kappa_0(\gamma) > 0 \), depending on \( \gamma \) and the other physical parameters in a semi-explicit way (see Theorem 3.19), such that

\[
(-\kappa_0(\gamma), \kappa_0(\gamma)) \subseteq \{ \kappa \in \mathbb{R} \mid (\gamma, \kappa) \in \mathfrak{P}^s \}.
\]

However, the estimates in Theorem 3.19 suggest that for each \( \kappa \in \mathbb{R} \setminus \{0\} \), the set \( \{ \gamma \in \mathbb{R} \setminus \{0\} \mid (\gamma, \kappa) \in \mathfrak{P}^s \} \) is bounded, and possibly empty for large \( |\kappa| \). We conjecture that this is indeed the case, but we do not have a complete proof here due to the complicated dependence the operator norm of \( L_{\sigma_{b+n}} \), defined in (3.65), on \( \gamma \).

Using solutions constructed in the flattened domain via Theorem 1.2, we can then produce solutions to the unflattened system (1.16). This leads us to the final result of this paper.
Theorem 1.3 (Proved later in Section 4.3). Let $\Sigma$ be given by (1.1), and suppose that $\mathbb{N} \ni s > \frac{n}{2}$. Suppose that either $\sigma > 0$ and $n \geq 2$ or else $\sigma = 0$ and $n = 2$. Let

$$\mathcal{U}^g \subset (\mathbb{R} \setminus \{0\}) \times \mathbb{R} \times H^{n+2}(\Sigma \times \mathbb{R}; \mathbb{R}^{n \times n}_{\text{sym}}) \times H^{n+\frac{1}{2}}(\Sigma; \mathbb{R}^{n \times n}_{\text{sym}}) \times H^{n+1}(\Sigma \times \mathbb{R}; \mathbb{R}) \times H^s(\Sigma; \mathbb{R}^n)$$

and $\mathcal{O}^g \subset \mathcal{X}^g$ be the same open sets as in Theorem 1.2. Then for each $(\gamma, \kappa, T, \tau, f, g) \in \mathcal{U}^g$, there exist

- a free surface function $\eta \in X^{s+\frac{1}{2}}(\Sigma; \mathbb{R}) \cap C^{(s-n/2)+1}(\Sigma; \mathbb{R})$ such that $\max_X |\eta| \leq b/2$, and $\mathcal{F}_\eta$ as defined in (1.17) a bi-Lipschitz homeomorphism and a $C^{(s-n/2)+3}$ diffeomorphism,
- a velocity field $v \in \mathcal{O}^{s+\frac{1}{2}}(\Omega_{b+\gamma; \mathbb{R}}^{n}; \mathbb{R}^n) \cap C^{(s-n/2)+2}(\Omega_{b+\gamma; \mathbb{R}}^{n})$,
- a pressure $q \in H^{n+1}(\Omega_{b+\gamma; \mathbb{R}}^{n}) \cap C^{(s-n/2)+1}(\Omega_{b+\gamma; \mathbb{R}}^{n})$,
- and constants $C, R > 0$.

such that the following hold:

1. $(v, q, \eta)$ are classical solutions to

$$\begin{align*}
\text{div } S(q, v) - \gamma e_1 \cdot \nabla v - \kappa \kappa x_n \partial_1 q e_1 + (v + W^1 + W^2) \cdot \nabla (v + W^1 + W^2) + (\nabla' \eta, 0) \\
- \kappa \kappa x_n \Delta' q e_1 - \kappa \kappa x_n \nabla \partial_1 q, \partial_1 \eta = f + L_{b+\gamma, f}, \\
\text{div } v - \kappa \kappa x_n \partial_1 \eta = 0, \\
(\kappa + \partial_n) + \kappa (\kappa + \partial_n) \partial_1 \eta = \tau' \mathcal{N}, \\
(Q, v) \mathcal{N} = [-\sigma \partial_q \mathcal{N} \mathcal{L} + T \partial_q + S b + \tau' T + \kappa (\kappa + \partial_n) (\kappa e_1 + (\nabla \eta, 0) + (\nabla' \eta, 0) \cdot e_1)] \mathcal{N}, \\
v = 0,
\end{align*}$$

on $\Omega_{b+\gamma}$, $\Omega_{b+\gamma}$, $\Omega_{b+\gamma}$, $\Omega_{b+\gamma}$, $\Sigma_0$.

(2) $(v \circ \mathcal{F}_\eta, q \circ \mathcal{F}_\eta, \eta) \in \mathcal{O}^g \subset \mathcal{X}^g$.

(3) If $(v_*, \kappa_*, T_*, \tau_*, f_*, g_*) \in \mathcal{U}^g$ satisfy

$$|\gamma - \gamma_*| + |\kappa - \kappa_*| + \|T - T_*\|_{H^{s+2}} + \|T - T_*\|_{H^{s+1/2}} + \|f - f_*\|_{H^{s+1}} + \|f - f_*\|_{H^{s}} < R,$$

then the corresponding solution triple $(v_*, q_*, \eta_*)$ satisfy

$$\|(v - \mathcal{F}_\eta, q - \mathcal{F}_\eta, \eta) - (v_* \circ \mathcal{F}_\eta, q_* \circ \mathcal{F}_\eta, \eta_*)\|_{\mathcal{X}} \leq C(|\gamma - \gamma_*| + |\kappa - \kappa_*| + \|T - T_*\|_{H^{s+2}} + \|T - T_*\|_{H^{s+1/2}} + \|f - f_*\|_{H^{s+1}} + \|f - f_*\|_{H^{s}}).$$

1.7. Discussion of techniques and plan of attack. As in [18], our strategy for producing solutions to (1.24) is to employ an implicit function theorem argument built on the linearization of (1.24),

$$\left\{\begin{array}{ll}
\text{div } S(p, u) - \gamma \partial_1 u - \kappa \kappa x_n \partial_1 \eta e_1 + (\nabla' \eta, 0) + s(x_n) \partial_1 u \\
+ s(x_n) u_e e_1 + \kappa \kappa x_n (\kappa e_1 - \kappa \kappa x_n \nabla' \eta e_1 - \kappa \kappa x_n \nabla' \eta, \partial_1 \eta) = f, & \text{in } \Omega \\
\text{div } u + \kappa \kappa x_n \partial_1 \eta = g, & \text{in } \Omega \\
u_n + (\kappa \kappa x_n^2 \partial_1 \eta, e_1 = h, & \text{on } \Sigma_0 \\
S(p, u)e_n + \sigma \Delta' \eta e_n = k, & \text{on } \Sigma_0 \\
u_0 = 0, & \text{on } \Sigma_0.
\end{array}\right.$$
for $\xi \in \hat{\Sigma}$, where $\hat{\Sigma}$ is the dual group of $\Sigma$ defined via (1.45), and $\psi, \rho_\gamma$ are defined in (3.35), (3.19) as

$$\psi(\xi) = \int_0^b \left( \hat{f}(\xi, x_n) \cdot \nabla(\xi, x_n, -\gamma) - \hat{g}(\xi, x_n)Q(\xi, x_n, -\gamma) \right) \, dx_n - \hat{k}(\xi) \cdot \nabla(\xi, b, -\gamma) + \hat{h}(\xi),$$

$$\rho_\gamma(\xi) = 2\pi i \gamma_1 + (1 + 4\pi^2|\xi|^2 \sigma) m(\xi, -\gamma).$$

(1.38)

Here $V, Q, m$, defined via (3.9), are symbols of special pseudodifferential operators. In order to solve the pseudodifferential equation for $\eta$, we need the precise asymptotics of $V, Q, m$ on $\hat{\Sigma}$. In particular, we note that when $\gamma \neq 0$ and $\Sigma$ is defined via (1.1), the asymptotics of $m$ completely determine the asymptotics of $\rho_\gamma$ on $\hat{\Sigma}$, which for $\gamma \neq 0$ are given in Lemma 3.6 as

$$\begin{align*}
|\rho_\gamma(\xi)|^2 &\asymp (\xi_1^2 + |\xi|^4) \mathbb{I}_{B_{\Sigma}(0, 1)}(\xi) + (1 + |\xi|^2) \mathbb{I}_{B_{\Sigma}(0, 1)^c}(\xi), & \sigma > 0 \\
|\rho_\gamma(\xi)|^2 &\asymp |\xi|^2 \mathbb{I}_{B_{\Sigma}(0, 1)}(\xi) + (1 + |\xi|^2) \mathbb{I}_{B_{\Sigma}(0, 1)^c}(\xi), & \sigma = 0, n = 2.
\end{align*}$$

(1.39)

With the asymptotics of $\rho_\gamma$ in hand, we define $\hat{\eta}$ via $\hat{\eta}(\xi) = \psi(\xi)/\rho(\xi)$ for $\xi \neq 0$, which by (1.39) would imply

$$\int_{\Sigma} \frac{\xi_1^2 + |\xi|^4}{|\xi|^2} |\hat{\eta}(\xi)|^2 \mathbb{I}_{B_{\Sigma}(0, 1)} + (1 + |\xi|^2)^{s+5/2} |\eta(\xi)|^2 \mathbb{I}_{B_{\Sigma}(0, 1)^c} \, d\xi \asymp \int_{\Sigma} \frac{|\psi(\xi)|^2}{|\xi|^2} \mathbb{I}_{B_{\Sigma}(0, 1)} + (1 + |\xi|^2)^{s+3/2} |\psi(\xi)|^2 \mathbb{I}_{B_{\Sigma}(0, 1)^c} \, d\xi.$$  

(1.40)

Using the asymptotics of $V, Q$ in Theorem 3.5, and the functional framework built on the specialized anisotropic space $X^s(\Sigma; \mathbb{R})$ that serves as the container space for the free surface function $\eta$ in Section 2.1, we can then utilize the equivalence (1.40) to recover $\eta$ via Fourier reconstruction. We note that by the first items of Lemma 2.2 and Theorem 2.6, the free surface function $\eta$ recovered through this process will be regular enough to be a classical function as opposed just being a tempered distribution, which is a crucial requirement for $\eta$ to be utilized in the subsequent nonlinear analysis.

It is also worth mentioning that in the physically relevant dimension $n = 3$, we deliberately chose to ignore the configuration $\Sigma = (L_1 \mathbb{T}) \times \mathbb{R}$, as in this case we would have $\xi \in B_2(0, r) \iff \xi = (0, \xi_2), |\xi_2| \leq r$ for $r = \min(1, L_1^{-1})$. By (1.38), this implies that $|\rho_\gamma(\xi)|^2 \mathbb{I}_{B_{\Sigma}(0, r)} \asymp |\xi|^4 \iff |\xi|^2 |\hat{\eta}(\xi)|^2 \mathbb{I}_{B_{\Sigma}(0, r)} \asymp \frac{|\psi(\xi)|^2}{|\xi|^2} \mathbb{I}_{B_{\Sigma}(0, r)}$. Unfortunately, in this scenario we only have $H^1$ control over the low frequency modes of $\eta$, and by Proposition 2.3 the corresponding anisotropic space $X^s((L_1 \mathbb{T} \times \mathbb{R})\mathbb{R})$ fails to be complete. If one were to consider the completion of this space, elements of the completion would be equivalence classes of tempered distributions modulo polynomials. Since we mandate elements of the container space for the free surface function $\eta$ to be classical functions, neither the space $X^s((L_1 \mathbb{T} \times \mathbb{R})\mathbb{R})$ nor its completion are suitable for the purposes of this paper. Practically speaking, this means that we cannot employ our techniques assuming a priori periodicity in the direction of incline. For the same reason, our framework also cannot produce stationary solutions in the case when we cannot employ our techniques assuming a priori periodicity in the direction of incline. For the same reason, our framework also cannot produce stationary solutions in the case when we cannot employ our techniques assuming a priori periodicity in the direction of incline. For the same reason, our framework also cannot produce stationary solutions in the case when we cannot employ our techniques assuming a priori periodicity in the direction of incline.

By solving for $\eta$ through the aforementioned approach, we can then solve for $(u, p)$ by means of (1.36) and the linear isomorphism associated to (1.37). Fortunately, this is possible as the results from the linear analysis in [18] continue to hold over $\Sigma$ defined in (1.1). This is mainly due to the fact that many results in [18] are proved by analyzing the low frequency behavior of functions in $H^s(\mathbb{R}^d; \mathbb{R})$ and $X^s(\mathbb{R}^d; \mathbb{R})$. By Remark A.6 and Remark 3.2, the analogous results for $H^s(\Sigma; \mathbb{R})$ and $X^s(\Sigma; \mathbb{R})$ can all be deduced from reducing a similar set of calculations over $\Sigma$ and $\hat{\Sigma}$ to the calculations in [18] over $\mathbb{R}^d$. As such, the results from [18] can be directly ported over with minimal modification of the proofs contained therein.

Though, as seen in the case $\Sigma = L_1 \mathbb{T} \times \mathbb{R}$, the generalized space $X^s(\Gamma; \mathbb{R})$ can fail to be complete depending on the set $\Gamma$ over which it is defined. This is an initial indication that some care needs to be taken in generalizing the space $X^s(\mathbb{R}^d; \mathbb{R})$ introduced in [18] over general sets $\Gamma$. We devote the Appendix to developing the necessary tools to carefully define the space $X^s(\gamma; \mathbb{R})$ over product domains $\Gamma$ in Section 2.1, and also identify the “good sets” $\Sigma$ for which the spaces $X^s(\Sigma; \mathbb{R})$ are compatible with our analysis, which led us to (1.1).

Next, we discuss the role of the perturbations (1.9), (1.15), and in particular the role of (1.14). First, by renormalizing the pressure via (1.8), the vertical gravitational force $-c_n$ shifts from the bulk to the dynamic boundary conditions, appearing as the term $\zeta \psi$ on the right hand side of the fourth equation in (1.10); this term later appears as $\eta N$ in the fourth equation of (1.12). By further modifying the pressure via (1.15), we eliminate the term $\eta N$ from (1.12), at the price of introducing $(\nabla \psi, 0)$ to the bulk in (1.16) and (1.24). The advantage of this formulation is that the pressure $q$ in Theorem 1.3 lives in a standard Sobolev space $H^s(\Omega; \mathbb{R})$, as opposed to the alternative formulation in [18], for which $q$ belongs to a specialized Sobolev space built on the anisotropic space $X^s(\Sigma; \mathbb{R})$.  

TRAVELING WAVE SOLUTIONS TO FREE BOUNDARY INCOMPRESSIBLE NAVIER-STOKES EQUATIONS 9
A key difference between the system \((1.10)\) and the analogous system in \([18]\) is that upon removing the background shear flow \((1.7)\), we are left with the term \(\kappa \eta_1 \delta e_1 \otimes e_n + e_n \otimes \delta e_1\) appearing in the fourth equation of \((1.12)\). Since \(\eta\) is expected to belong to a specialized Sobolev space \(X^s(\Sigma; \mathbb{R})\) that does not coincide with the standard Sobolev space \(H^s(\Sigma; \mathbb{R})\) in general, attacking the problem at this level would require one to build specialized spaces for the data tuple \((f, g, h, k)\), and also prove that the associated linear maps remain to be isomorphisms in this modified functional framework. By introducing an additional perturbation \((1.14)\), we were able to replace this term with functions in \(s > n/2\) thanks to the second item of Theorem 2.6, as the non-trivial terms are either derivatives of functions in \(X^s(\Sigma; \mathbb{R})\) or products of functions in \(X^s(\Sigma; \mathbb{R})\) and derivatives of functions in \(X^s(\Sigma; \mathbb{R})\). This approach allows us to directly employ the functional framework from \([18]\) at the price of introducing worse nonlinearities in \((1.16)\) and \((1.24)\).

To construct solutions to \((1.16)\), we use the associated linear isomorphism in conjunction with the implicit function theorem around the trivial solution. In order to invoke the implicit function theorem, the nonlinear maps associated to \((1.24)\) first need to be well-defined on the same spaces used in the linear analysis. This requires some analysis in Section 4.1 to understand the mapping properties of functions in the specialized Sobolev space \(X^s(\Sigma; \mathbb{R})\), and also of the mean-curvature operator \(\mathcal{H}\) defined in \((1.2)\).

In addition, we also need a divergence-trace type compatibility condition \((3.22)\) built into the container space \(\mathcal{Y}^s\) defined in \((3.23)\) to hold; in Theorem 4.2, we show that the compatibility condition requires

\[
\left[ -s(\eta + b)\partial_1 \eta - \kappa (\eta + b)\partial_1 \eta - \int_0^1 J(x', x_n) \kappa \left( x_n + \eta(x') \frac{x_n}{b} \right) \partial_1 \eta(\cdot) \, dx_n \right]_{\mathcal{H}^{-1}} = \left[ -\kappa (b^2 \partial_1 \eta + b \partial_1 \eta^2 + \frac{1}{3} \partial_3 \eta^3) \right]_{\mathcal{H}^{-1}} < \infty, \tag{1.41}
\]

where the \(\mathcal{H}^{-1}\) seminorm is defined in \((3.2)\).

A major obstacle in proving \((1.41)\) is that low-frequency control of powers of \(\eta\) is not immediately evident from the inclusion \(\eta \in X^s(\Sigma; \mathbb{R})\). One could in principle attempt to obtain control of powers of \(\eta\) by way of Young’s convolution inequality applied on the Fourier side. However, in the model case \(\Sigma = \mathbb{R}^d\) with \(s > d/2\), this only leads to the inclusion \(\eta^k \in H^k(\mathbb{R}^d; \mathbb{R})\) for \(k \geq 2 + [4/(2d - 2)], d \geq 2\). Unfortunately, this means that in the physically relevant case \(n = 3\) (so that \(d = n - 1 = 2\)), this elementary argument only provides control over quartic or higher powers of \(\eta\). However, by the fourth item of Theorem 2.6, we know that if \(\zeta \in X^s(\Sigma; \mathbb{R})\), then \(\partial_1 \zeta \in \mathcal{H}^{-1}(\Sigma; \mathbb{R})\). Thus, a viable strategy for proving \((1.41)\) is to show that \(X^s(\Sigma; \mathbb{R})\) is an algebra. By Lemma 2.2, there are configurations of \(\Sigma\) for which \(X^s(\Sigma; \mathbb{R}) = H^s(\Sigma; \mathbb{R})\), in which case we know \(X^s(\Sigma; \mathbb{R})\) is an algebra for \(s > \text{dim } \Sigma/2\). In general, though, we only know that \(H^s(\Sigma; \mathbb{R}) \hookrightarrow X^s(\Sigma; \mathbb{R})\), so further analysis is required to show that \(X^s(\Sigma; \mathbb{R})\) is an algebra. Fortunately, we are able to establish this in Theorem 1.1, which is proved later in Section 2.2.

As the linearization of \((1.24)\) depends on \(\kappa\), we also need to identify a parameter regime for \(\kappa\) for which the associated linear map remain an isomorphism. To that end, in Section 3.6 we study the map

\[
L_{\kappa, \sigma}(u, p, \eta) = (\text{div } S(p, u) - \gamma \kappa \eta_n \partial_1 \eta e_1 - \gamma \partial_1 u + (\nabla' \eta, 0) + s(x_n) \partial_1 u + s'(x_n) u_n e_1 + \kappa x_n s(x_n) \partial_1 \eta e_1

- \kappa x_n \Delta' \eta e_1 - \kappa (x_n \partial_1 \nabla' \eta, \partial_1 \eta) \delta u + \kappa x_n \partial_1 \eta u_n |\Sigma_b + (\gamma - \frac{\kappa \delta^2}{2}) \partial_1 \eta s(p, u) e_n |\Sigma_b + \sigma \Delta' \eta e_n) \tag{1.42}
\]

induced by the linearization associated to \((1.24)\), and in Theorem 3.19 we show that via a perturbative argument around \(\kappa = 0\) that for fixed \(\gamma\) and other physical parameters, there exists a \(\kappa_0 > 0\) for which \(L_{\kappa, \sigma}\) is an isomorphism over appropriate spaces for all \(\kappa \in (-\kappa_0, \kappa_0)\). Synthesizing the aforementioned results and employing our strategy of invoking the implicit function theorem leads us to the solvability of \((1.24)\) and \((1.16)\).

Finally, we discuss the strategy for producing solutions to the unflattened system \((1.16)\) using solutions constructed in the flattened domain via Theorem 1.2. To that end, we use the free surface function \(\eta\) to build the flattening map and its inverse defined in \((1.17)\) and \((1.18)\) to undo the reformulation outlined in Section 1.5. This requires some results on the regularity of these maps, which is recorded in Section 4.1. Fortunately, the same analysis can be adapted from \([18]\) with minimal modification.

1.8. Notational conventions and outline of article. In this subsection we discuss the notational conventions adopted throughout this paper. In this paper, \(\mathbb{N}\) denotes the natural numbers including 0. We always use \(n \geq 1\) to denote the dimension of the flattened fluid domain \(\Omega = \Sigma \times \mathbb{R}\). For \(d \geq 1\), we will consider spaces defined over
domains of the form
\[ \Gamma = \prod_{i=1}^{d} \Gamma_i, \quad \Gamma_i = \mathbb{R} \text{ or } L_i \mathbb{T}, \quad L_i > 0, \]  
where \( \Gamma \) is endowed with the natural group, topological, and smooth structures. In fact, the topology on \( \Gamma \) is metrizable, and in this paper we equip \( \Gamma \) with the metric
\[ d_\Gamma(x, y) = \left( \sum_{i=1}^{d} d_i(x_i, y_i)^2 \right)^{1/2} \]
for \( d_i(x_i, y_i) = \begin{cases} |x_i - y_i|, & \text{if } i \in R_G, \\ \inf\{|r - s| : r \in [x_i], s \in [y_i]|, & \text{if } i \in T_G. \end{cases} \]

We write \( \hat{\Gamma} \) to denote the dual group associated to \( \Gamma \), defined via
\[ \hat{\Gamma} = \prod_{i=1}^{d} \hat{\Gamma}_i, \quad \hat{\Gamma}_i = \begin{cases} \mathbb{R}, & \Gamma_i = \mathbb{R} \\ L_i^{-1} \mathbb{Z}, & \Gamma_i = L_i \mathbb{T}, \end{cases} \]
where \( \hat{\Gamma} \) is also endowed with the obvious group, topological, and smooth structures. We also endow \( \hat{\Gamma} \) with the metric induced by inclusion in \( \mathbb{R}^d \):
\[ d_{\hat{\Gamma}}(x, y) = \sum_{i=1}^{d} (|x_i - y_i|^2)^{1/2}. \]

If \( X \) is a metric space, we write \( B_X(x, r) = \{ y \in X : d_X(x, y) < r \} \) for the open ball.

For \( X \in \{ \Gamma, \hat{\Gamma} \} \), we write \( \mathcal{S}(X; \mathbb{C}) \) to denote the Schwartz class of complex valued functions over \( X \) and \( \mathcal{S}'(X; \mathbb{C}) \) to denote the space of complex valued tempered distributions over \( X \); a detailed treatment of how to define these space can be found in Appendix A.2. For \( f \in \mathcal{S}(X; \mathbb{C}) \) or \( f \in \mathcal{S}'(X; \mathbb{C}) \), we denote its unitary Fourier and inverse Fourier transforms by
\[ \hat{f}(\xi) = \mathcal{F}_X(f)(\xi), \quad \check{f}(\xi) = \mathcal{F}_X^{-1}(f)(\xi), \]
where \( \mathcal{F}_X \) is defined in (A.19). Sometimes we will also write
\[ \mathcal{F}[f](\xi) = \mathcal{F}_X(f)(\xi), \quad \mathcal{F}^{-1}[f](\xi) = \mathcal{F}_X^{-1}(f)(\xi). \]

For \( \mathbb{R} \ni s > 1/2 \) and a Lipschitz \( \zeta : \Sigma \to \mathbb{R} \) satisfying \( \inf \zeta > 0 \), we define
\[ _0H^s(\Omega; \mathbb{R}^n) = \{ u \in H^s(\Omega; \mathbb{R}^n) : u = 0 \text{ on } \Sigma_0 \}, \]
where the equality is taken in the trace sense. In the case when \( s = 1 \) and \( \zeta = b \), we endow \( _0H^1(\Omega_b; \mathbb{R}^n) \) with the inner product
\[ (u, v)_{_0H^1} = \frac{1}{2} \int \mathbb{D} u : \mathbb{D} v \, dx, \]
which generates the same topology as the standard \( H^1 \)-norm by Korn’s inequality (see Lemma 2.7 of [3]). For \( k \in \mathbb{N} \), a real Banach space \( V \), and a nonempty open set \( U \subseteq \Gamma \) we define the space
\[ C^k(U; V) = \{ f : U \to V : Lf \in C^k(\mathbb{R}^d; V), \quad Lf(x) = f([x]) \}. \]

We define
\[ C^k_b(U; V) = \{ f : U \to V : f \in C^k(U; V) \text{ and } \| f \|_{C^k_b} < \infty \}, \]
where
\[ \| f \|_{C^k_b} = \sum_{\alpha \in \mathbb{N}^d, |\alpha| \leq k} \sup_{x \in U} \| \partial^\alpha f(x) \|_V. \]

We also define the space \( C^k_0(\Gamma; V) \subset C^k_b(\Gamma; V) \) to be the closed subspace
\[ C^k_0(\Gamma; V) = \{ f \in C^k_b(\Gamma; V) : \lim_{|x_{R_\Gamma}| \to \infty} \partial^\alpha f(x) = 0 \text{ for all } \alpha \in \mathbb{N}^d \text{ such that } |\alpha| \leq k \}, \]
where \( x_{R_\Gamma} \in \Gamma \) is defined via (A.2). Lastly, we use \( L^0(\Gamma; [0, \infty]) \) to denote the set of non-negative measurable functions over \( \Gamma \).

We conclude this section by giving outline of the article. In Section 2, we introduce the anisotropic Sobolev space \( X^s(\Sigma; \mathbb{R}) \) and characterize the space based on the underlying product structure of the domain \( \Gamma \). We then state its essential properties, and in Section 2.2 we prove that for \( \Sigma \) defined via (1.1), \( X^s(\Sigma; \mathbb{R}) \) is an algebra in the supercritical regime \( s > \dim \Sigma/2 \).
In Section 3, we first record the isomorphism associated to the overdetermined system (1.37) and the asymptotics of the special pseudodifferential symbols used in the construction of the free surface function $\eta$. This allows us to prove the isomorphism associated to the linear system (1.36). We then establish the parameter regime for $\kappa$ for which the flattened system (1.24) induces an isomorphism.

In Section 4, we first record some key mapping properties of the anisotropic space $X^s(\Sigma; \mathbb{R})$ and various nonlinear maps used in the subsequent analysis. Using these preliminary results, we show that the maps induced by (1.24) are well-defined and smooth, and we use this in conjunction with the implicit function theorem to produce solutions to (1.24). We conclude the paper by using the solutions from (1.24) to produce solutions to (1.16).

2. The anisotropic space $X^s(\Gamma; \mathbb{F})$

In this section we aim to generalize the anisotropic Sobolev space $X^s(\mathbb{R}^d; \mathbb{R})$ introduced in Section 5 of [18]. First, in order to define such a space on more general product domains, we need to define a suitable notion of tempered distributions on $\Gamma$ and its Pontryagin dual $\hat{\Gamma}$; a detailed treatment of this can be found in Appendix A.2.

First, in order to define such a space on more general product domains, we need to define a suitable notion of the class of tempered distributions $\mathcal{S}(\Gamma; \mathbb{F})$ with the norm associated to $\mu$. We note that defining $\mu$ over $\hat{\Gamma}$ is well-defined and smooth, and we use this in conjunction with the implicit function theorem to produce solutions to (1.24).

The definitions of the class of tempered distributions $\mathcal{S}(\Gamma; \mathbb{C})$ on $\Gamma$ and the Fourier transform on $\mathcal{S}'(\Gamma; \mathbb{C})$ are contained in Section A.2 of the appendix.

First, we show that if a permutation $g \in S_d$ leaves the first factor of $\Gamma$ fixed, then the restriction of the induced map $\sigma_g$ on $X^s(\Gamma; \mathbb{F})$ is an isometric isomorphism between $X^s(\Gamma; \mathbb{F})$ and $X^s(g\Gamma; \mathbb{F})$. The reordered Cartesian product $g\Gamma$ and the map $\sigma_g$ are defined in Definition A.3.
Lemma 2.1. Let $\Gamma$ be defined as in (1.43) and let $S_d$ denote the symmetric group of permutations. Suppose $g \in S_d$ satisfies $g(1) = 1$, and let $g^\Gamma$ be the reordered Cartesian product defined as in the first item of Definition A.3. \( \mathcal{P}_g : \Gamma \to g^\Gamma \) be the reordering map defined as in the second item of Definition A.3. Then the map \( \sigma_g : X^s(\Gamma; \mathbb{F}) \to X^s(\Gamma; \mathbb{F}) \) defined via \( \sigma_g(f) = f(P_{\mathcal{P}_g}^{-1}) \) is an isometric isomorphism, with its inverse given by the map \( \sigma_g^{-1} : X^s(\mathcal{P}_g; \mathbb{F}) \to X^s(\Gamma; \mathbb{F}) \) defined via \( \sigma_g^{-1}(f) = f(P_g) \).

Proof. We note that \( \langle \cdot \rangle = \langle \mathcal{P}_g \rangle \), and if \( g(1) = 1 \), then \( \mu(\cdot) = \mu(\mathcal{P}_g) \). As a consequence, the Fourier multiplier \( \mu^2(\cdot)\langle \cdot \rangle^{2(s-1)} \) is invariant under \( \mathcal{P}_g \), therefore the desired conclusion follows from Lemma A.11.

Next, we characterize the anisotropic space \( X^s(\Gamma; \mathbb{F}) \) in relation to the set \( \Gamma_r \) defined in the first item of Definition A.2.

Lemma 2.2. Let $\Gamma$ be defined as in (1.43) with dim $\Gamma = d \geq 1$ and let $\Gamma_r$ be defined as in (A.2). Let $\Gamma_\ast$, $\Gamma_r$ be the canonical reordering of $\Gamma$ and $\Gamma_r$ defined in (A.8), and let $\mathcal{P}_\Gamma : \Gamma \to \Gamma_\ast$, $\mathcal{P}_r : \Gamma \to \Gamma_r$ be the canonical reordering maps defined in (A.10). Then the following hold.

1. If $R \Gamma = \emptyset$ or $R \Gamma = \{1\}$, then $X^s(\Gamma; \mathbb{F}) = H^s(\Gamma; \mathbb{F})$ and $\| \cdot \|_{X^s}$ and $\| \cdot \|_{H^s}$ are equivalent norms. In particular, this implies that if $s > d/2$, then $X^s(\Gamma; \mathbb{F})$ is an algebra.

2. If $1 \in R \Gamma$ and $|R \Gamma| \geq 2$, then $X^s(\Gamma; \mathbb{F})$ is not closed under rotations in the sense that for any orthogonal matrix $Q \in O(|R \Gamma|)$ such that $|Qe_1 \cdot e_1| < 1$, there exists a function $f \in X^s(\Gamma; \mathbb{F}) \setminus L^2(\Gamma; \mathbb{F})$ such that $f(P_{\Gamma_r}^{-1}(Q \oplus I)P_{\Gamma_r}) \notin X^s(\Gamma; \mathbb{F})$, where $Q \oplus I : \Gamma_\ast \oplus \Gamma_r \to \Gamma_\ast \oplus \Gamma_r$ maps $(x, y)$ to $(Qx, y)$. In particular, this implies that $H^s(\Gamma; \mathbb{F}) \subseteq X^s(\Gamma; \mathbb{F})$.

3. If $1 \notin R \Gamma$ and $R \Gamma \neq \emptyset$, then $H^s(\Gamma; \mathbb{F}) \subseteq X^s(\Gamma; \mathbb{F})$. Furthermore, for all $f \in X^s(\Gamma; \mathbb{F})$ we have

$$
\|f\|_{X^s} \leq \int_\Gamma |\xi|^2 \langle \xi \rangle^{2(s-1)} |\hat{f}(\xi)|^2 \, d\xi. \tag{2.7}
$$

Proof. We note that since $\mu(\cdot) \leq 1 + |\xi|$, we always have $H^s(\Gamma; \mathbb{R}) \subseteq X^s(\Gamma; \mathbb{R})$. To prove the first item, we note that if $R \Gamma = \emptyset$, then we are in the purely toroidal case when $\Gamma = \prod_{i=1}^d L_i R$. Here, the only low frequency mode $\xi \in B(0, r) = 0$. By the definition of $\mu$ at $\xi = 0$ in (2.4), we have $\mu(0) = 1 = (0)^{2s}$. So by (2.2) we find that for all $\xi \in B(0, r)$, $\mu(\xi) \sim |\xi|^{2s}$ and the desired conclusion follows. If $R \Gamma = \{1\}$, then $\Sigma = \mathbb{R} \times \prod_{i=1}^d L_i R$. In this case, $\xi \in B(0, r) \iff \xi = (\xi_1, 0)$ for $|\xi_1| < r$. Then for $\xi \in B(0, r)$, we have $\mu(\xi) = 1 + |\xi_1|$ and $\langle \xi \rangle = (1 + |\xi_1|^2)^{1/2}$, which implies that $\mu(\xi) \sim |\xi_1|^2 B(0, r) + |\xi_1|^{2s} B(0, r) \sim |\xi|^{2s}$ for all $\xi \in \Gamma_\ast$, and the desired conclusion follows.

Next we proceed to prove the second item. Let the set $T \Gamma$ be defined as in (A.3). The case when $T \Gamma = \emptyset$ follows from the third item of Proposition 5.2 in [18], so we assume here that $T \Gamma \neq \emptyset$.

Since $|R \Gamma| \geq 2$, by the third item of Proposition 5.2 in [18], for every $Q \in O(|R \Gamma|)$ such that $|Qe_1 \cdot e_1| < 1$ there exists a function $G \in X^s(\mathbb{R}^{R \Gamma}; \mathbb{F}) \cap C^\infty_0(\mathbb{R}^{R \Gamma}; \mathbb{F})) \setminus L^2(\mathbb{R}^{R \Gamma}; \mathbb{F})$ such that $G(Q) \notin X^s(\mathbb{R}^{R \Gamma}; \mathbb{F})$. Let $G$ be such a function, $P_{\Gamma}$ be the canonical reordering map defined via (A.10), $\Gamma_r$ be the set defined via (A.4), and $\pi_r$ be the surjective map on $\hat{\Gamma}$ defined via (A.14). We then consider the measurable function $F : \hat{\Gamma} \to \mathbb{C}$ defined via

$$
F(\xi) = \begin{cases} 
\hat{G}(\pi_r \xi), & \xi \in \Gamma_r, \\
0, & \text{otherwise,}
\end{cases} \tag{2.8}
$$

where $\Gamma_r$ is defined as in (A.4). We note that by Lemma A.6, $\pi_r$ is an isometric measure-preserving group isomorphism between $\hat{\Gamma}_r$ and $\hat{\mathbb{R}}^{\Gamma_r} = \hat{\mathbb{R}}^{|R \Gamma|}$. Then by Fubini's theorem, it follows that

$$
\int_{\Gamma} |F(\xi)| \, d\xi = \int_{\Gamma_r} |\hat{G}(\pi_r \xi)| \, d\xi = \int_{\mathbb{R}^{|R \Gamma|}} |\hat{G}(\omega)| \, d\omega < \infty, \quad \int_{\Gamma} |F(\xi)|^2 \, d\xi = \int_{\mathbb{R}^{|R \Gamma|}} |\hat{G}(\omega)|^2 \, d\omega = \infty. \tag{2.9}
$$

We also note that the Fourier multiplier $\mu^2(\cdot)\langle \cdot \rangle^{2(s-1)}$ is invariant under $P_{\Gamma}$, and since $\pi_r = \pi_r \circ P_{\Gamma}$, by the same calculations in (A.48) and by Lemma A.6 we have

$$
\int_{\Gamma_r} \frac{1}{2} \mu^2(\xi)\langle \xi \rangle^{2(s-1)} |F(\xi)|^2 \, d\xi = \int_{\Gamma_r} \mu^2(\xi)\langle \xi \rangle^{2(s-1)} |\hat{G}(\pi_r \xi)|^2 \, d\xi = \int_{\mathbb{R}^{|R \Gamma|}} \mu(\omega)^2(\omega)^{2(s-1)} |\hat{G}(\omega)|^2 \, d\omega < \infty. \tag{2.10}
$$

Hence, the function $f := \hat{F} \in X^s(\Gamma; \mathbb{F}) \setminus L^2(\Gamma; \mathbb{F})$. In particular, $f \notin H^s(\Gamma; \mathbb{F})$. On the other hand, we have

$$
\mathcal{F}^s[\langle f(P_{\Gamma_r}^{-1}(Q \oplus I)P_{\Gamma_r}) \rangle](\xi) = \int_{\Gamma} f(P_{\Gamma_r}^{-1}(Q \oplus I)P_{\Gamma_r} x) e^{-2\pi i \xi \cdot x} \, dx = \int_{\Gamma_r} f(P_{\Gamma_r}^{-1}(Q \oplus I)P_{\Gamma_r} x) e^{-2\pi i \xi \cdot x} \, dx = \int_{\Gamma} f(x) e^{-2\pi i \xi \cdot (Q \oplus I)P_{\Gamma_r} x} \, dx. \tag{2.11}
$$
Then by Lemma A.6 and Fubini’s theorem,
\[
\|f(\mathcal{P}_T^{-1}(Q \oplus I)\mathcal{P}_T^{-1})\|^2_{L^2(\Gamma)} = \int_{\Gamma} \mu(\xi)\|\xi\|^{2(\text{s}-1)}|f(\mathcal{P}_T^{-1}(Q \oplus I)\mathcal{P}_T\xi)|^2 d\xi = \int_{\Gamma} \mu(\xi)\|\xi\|^{2(\text{s}-1)}|\hat{G}(\pi_\text{T}^{-1}(Q \oplus I)\mathcal{P}_T\xi)|^2 d\xi = \int_{\Gamma} \mu(\xi)\|\xi\|^{2(\text{s}-1)}|\hat{G}(\pi_\text{T}^{-1}(Q \oplus I)\mathcal{P}_T\xi)|^2 d\xi
\]
This proves the second item.

To prove the last item, we note that if \(1 \notin R_T\) and \(R_T \neq \emptyset\), then \(\Gamma_1 = L_1\mathbb{I}\) and \(\Gamma \neq \bigcup_{i=1}^d L_i\mathbb{I}\). In this setting, \(\xi \in B_\Gamma(0,r) \iff \xi \in \Gamma\) as \(\xi_T = 0\) for the low frequencies and \(B_\Gamma(0,r) \setminus \{0\} \neq \emptyset\). In particular, for \(\xi \in B_\Gamma(0,r)\) we have \(\xi_1 = 0\), which implies that \(\mu(\xi) = |\xi_{R_T}| = |\xi|\) for \(\xi \in B_\Gamma(0,r) \setminus \{0\}\), so by (2.2), we arrive at (2.7). By (2.7), it also follows that \(H^s(\Gamma;\mathbb{F}) \subseteq X^s(\Gamma;\mathbb{F})\).

Recall that our ultimate aim in introducing \(X^s(\Gamma;\mathbb{F})\) is to use it as the container space for the free surface function \(\eta\) in our study of the traveling wave problem. The upshot of Lemma 2.2 is that the precise structure of the space \(X^s(\Gamma;\mathbb{F})\) is heavily dependent on the form of the domain \(\Gamma\), and in particular, the properties of the set \(R_T\). In the first case considered in the lemma, \(X^s(\Gamma;\mathbb{F})\) is the standard Sobolev space \(H^s(\Gamma;\mathbb{F})\), and therefore we may employ standard Sobolev tools in our subsequent analysis. In the second case, even though \(X^s(\Gamma;\mathbb{F})\) is not a standard Sobolev space, we will be able to prove that it enjoys many of the same properties as \(H^s(\Gamma;\mathbb{F})\), however, in the third case, which includes the physically relevant case when \(\Gamma = L_1\mathbb{I} \times \mathbb{R}\), the space \(X^s(\Gamma;\mathbb{F})\) is unfortunately unusable for our subsequent analysis due to a failure of completeness. To justify the last claim, we prove the following proposition.

**Proposition 2.3.** Let \(\Gamma\) be defined as in (1.43) and let \(R_T\) be defined as in (A.2). If \(1 \notin R_T\) and \(R_T \neq \emptyset\), then \(X^s(\Gamma;\mathbb{F})\) is complete if and only if \(|R_T| > 2\).

**Proof.** The proof is a modification of Proposition 1.34 in [2], which characterizes when the homogeneous Sobolev space \(H^s(\mathbb{R}^d)\) is complete. We define the measure \(\mu_s := \|\xi\|^2(\text{s}-1) d\xi\), and denote by \(L^s_\sharp(\hat{\Gamma};\mathbb{C})\) the complex valued functions in \(L^2(\hat{\Gamma};\mu_s)\).

First assume \(|R_T| > 2\), and suppose \(\{u_n\}_n\) is a Cauchy sequence in \(X^s(\Gamma;\mathbb{F})\). Then by (2.7), \(\{u_n\}_n\) is a Cauchy sequence in \(L^s_\sharp(\hat{\Gamma};\mathbb{C})\), and hence there exists a function \(f \in L^s_\sharp(\hat{\Gamma};\mathbb{C})\) such that \(u_n \to f\) in \(L^s_\sharp(\hat{\Gamma};\mathbb{C})\). Recall that by (A.13), we have \(\xi \in B_\Gamma(0,r) \iff \xi \in \hat{\Gamma}\), where \(\hat{\Gamma}\) is defined in (A.4). Using Lemma A.6 and the assumption \(|R_T| > 2\), we have
\[
\int_{B_\Gamma(0,r)} |f(\xi)| d\xi \leq \left(\int_{B_\Gamma(0,r)} |\xi|^2|f(\xi)|^2 d\xi\right)^{1/2} \left(\int_{B_{R_T}(0,r)} |\omega|^2 d\omega\right)^{1/2} < \infty,
\]
where in the last equality we have used the isometric measure-preserving isomorphism \(\pi_\text{T}\) between \(B_\Gamma(0,r)\) and \(B_{R_T}(0,r)\]. Since \(\Gamma_{B_\Gamma(0,r)} f\) belongs to \(H^s(\hat{\Gamma};\mathbb{C})\), with (2.13) we can infer that \(f\) defines a tempered distribution.

Then \(u := f \in X^s(\Gamma;\mathbb{F})\) and \(u_n \to u\) in \(X^s(\Gamma;\mathbb{F})\). This shows that \(X^s(\Gamma;\mathbb{F})\) is complete.

Now assume \(|R_T| \leq 2\) and suppose for the sake of contradiction that \(X^s(\Gamma;\mathbb{F})\) is complete with respect to the \(X^s\) norm defined via (2.6). Consider a norm on \(X^s(\hat{\Gamma};\mathbb{F})\) given by \(\|u\|_* = \|u\|_{X^s} + \|\hat{u}\|_{L^1(B_\Gamma(0,r))}\); this norm is well-defined since the definition of \(X^s(\Gamma;\mathbb{F})\) requires \(\hat{u} \in L^1_{\text{loc}}(\hat{\Gamma};\mathbb{C})\) for any function \(u \in X^s(\Gamma;\mathbb{F})\). We then claim that \(X^s(\Gamma;\mathbb{F})\) is also complete endowed with the \(\|\cdot\|_*\) norm. Indeed, suppose \(\{u_n\}_n\) is a Cauchy sequence in \((X^s(\Gamma;\mathbb{F}),\|\cdot\|_*\))\), then \(\{u_n\}_n\) is a Cauchy sequence in \((X^s(\Gamma;\mathbb{F}),\|\cdot\|_{X^s}\)) and \(\{\hat{u}_n\}_n\) is a Cauchy sequence in \(L^1(B_\Gamma(0,r);\mathbb{C})\). By the assumed completeness of \((X^s(\Gamma;\mathbb{F}),\|\cdot\|_{X^s}\))\), there exists a function \(u \in (X^s(\Gamma;\mathbb{F}),\|\cdot\|_{X^s}\))\) for which \(u_n \to u\) in \((X^s(\Gamma;\mathbb{F}),\|\cdot\|_{X^s}\))\). Similarly, there exists a function \(g \in L^1(B_{R_T}(0,r);\mathbb{C})\) for which \(\hat{u}_n \to g\) in \(L^1(B_{R_T}(0,r);\mathbb{C})\). Clearly, \(g = \hat{u}\) a.e. in \(L^1(B_{R_T}(0,r);\mathbb{C})\), therefore we can conclude that \(u_n \to u\) in \((X^s(\hat{\Gamma};\mathbb{F}),\|\cdot\|_*\))\).
This completes the proof of the claim.

We now know that \(X^s(\hat{\Gamma};\mathbb{F})\) is complete with respect to both \(\|\cdot\|_{X^s}\) and \(\|\cdot\|_*\). The identity map \(I : (X^s(\hat{\Gamma};\mathbb{F}),\|\cdot\|_*\) \to (X^s(\hat{\Gamma};\mathbb{F}),\|\cdot\|_{X^s}\))\) is trivially continuous, thus we can invoke the bounded inverse theorem to deduce the existence of a universal constant \(C > 0\) such that \(\|u\|_* \leq C\|u\|_{X^s}\) for all functions \(u \in X^s(\Gamma;\mathbb{F})\). In turn, this implies that
\[
\|\hat{u}\|_{L^1(B_{R_T}(0,r))} \leq C\|u\|_{X^s}\text{ for all }u \in X^s(\Gamma;\mathbb{F}).
\]

To derive a contradiction, we construct an explicit function \(f \in X^s(\Gamma;\mathbb{F})\) for which (2.14) is violated. For any \(\xi \in \hat{\Gamma}\), we adopt the convention in (A.1) and write \(\xi = \xi_{R_T} + \xi_{T_R}\). Let \(C = \{\xi \in \hat{\Gamma} : |\xi_{R_T}| = 0\text{ and }r/2 < |\xi_{T_R}| < r\},\)
which in particular implies that \(2^{-i}C \cap 2^{-j}C = \emptyset\) for any \(i, j \in \mathbb{N}\) such that \(i \neq j\). Now for every \(n \geq 1\) we then consider \(G_n \in L^1(\tilde{\Gamma}; \mathbb{C})\) defined via
\[
G_n(\xi) = \sum_{q=1}^{n} \frac{2q^{(1+|R_q|)/2}}{q} 1_{2^{-q}C}(\xi).
\]  
(2.15)

We note in particular that \(\overline{G_n(\xi)} = G_n(-\xi)\) and \(\text{supp} G_n \subseteq B_{\Gamma}(0, r) \subset \tilde{\Gamma}\), and so we may define the smooth and bounded function \(g_n : \Gamma_* \to \mathbb{R}\) via \(g_n = \mathcal{F}^{-1}[G_n]\). Let \(\mathcal{H}\) denote the Hausdorff measure over \(\tilde{\Gamma}\). Now we may readily calculate
\[
\|\hat{g}_n\|_{L^1(B_{\Gamma}(0, r))} = \sum_{q=1}^{n} \frac{2q^{(1+|R_q|)/2}}{q} 2^{-q|R_q|} \mathcal{H}|R_q|(\mathcal{C}) = \mathcal{H}|R_q|(\mathcal{C}) \sum_{q=1}^{n} \frac{2q^{(1+|R_q|)/2}}{q} \to \infty \text{ as } n \to \infty,
\]  
(2.16)
since \(|R_q| \leq 2\). On the other hand, for any \(n \geq 1\) we have
\[
\|g_n\|_{X^s}^2 = \int_{B_{\Gamma}(0, r)} |\xi|^2 |G_n(\xi)|^2 \, d\xi = \sum_{q=1}^{n} \int_{2^{-q}C} |\omega|^2 \left(\frac{2q^{(1+|R_q|)/2}}{q}\right)^2 \, d\omega
\]
\[
\quad \leq \mathcal{H}|R_q|(\mathcal{C}) \sum_{q=1}^{n} 2^{-2q}q^{-1+|R_q|/2} \left(\frac{2q^{(1+|R_q|)/2}}{q}\right)^2 = \mathcal{H}|R_q|(\mathcal{C}) \sum_{q=1}^{n} \frac{1}{q^2}. \tag{2.17}
\]
Therefore, \(\sup_{n \geq 1} \|g_n\|_{X^s} < \infty\) and by (2.14) and (2.16), we arrive at a contradiction. Thus, \(X^s(\Gamma; \mathbb{F})\) cannot be complete for \(|R_q| \leq 2\).
\[\Box\]

We now proceed to study the space \(X^s(\Gamma; \mathbb{F})\) in the second case of Lemma 2.2. We begin by stating a preliminary result.

**Lemma 2.4.** Let \(\Gamma\) be defined as in (1.43) and \(r\) be defined as in (A.13). Suppose \(\mathbb{R} \ni s \geq 0, 1 \leq R_q, \text{ and } |R_q| \geq 2\). The following hold.

1. For \(\mu : \tilde{\Gamma} \to \mathbb{R}\) defined in (2.4), we have
\[
\int_{B_{\Gamma}(0, r)} \frac{1}{\mu^2(\xi)} \, d\xi < \infty. \tag{2.18}
\]

2. For \(f \in X^s(\Gamma; \mathbb{F})\) we have the estimate
\[
\int_{B_{\Gamma}(0, r)} |\hat{f}(\xi)| \, d\xi + \left(\int_{\Gamma \setminus B_{\Gamma}(0, r)} (1 + |\xi|^2)^{s/2} |\hat{f}(\xi)|^2 \, d\xi\right)^{1/2} \lesssim_{d, s} \|f\|_{X^s}. \tag{2.19}
\]

In particular, if \(s > d/2\) then
\[
\|\hat{f}\|_{L^1(\tilde{\Gamma})} \lesssim_{d, s} \|f\|_{X^s}. \tag{2.20}
\]

**Proof.** For the first item, we note that by Lemma A.6,
\[
\int_{B_{\Gamma}(0, r)} \frac{1}{\mu^2(\xi)} \, d\xi \lesssim \int_{B_{\Gamma}(0, r)} \frac{|\xi|^2}{\xi^2 + |\xi|^4} \, d\xi = \int_{B_{\Gamma}(0, r)} \frac{|\xi R_q|^2}{\xi^2 + |\xi R_q|^4} \, d\xi = \frac{1}{\mathcal{H}|R_q|(\mathcal{C})} \int_{B_{\tilde{\Gamma}|R_q|(0, r)}} \frac{|\omega|^2}{\omega^2 + |\omega|^4} \, d\omega, \tag{2.21}
\]
where in the last equality we used the isometric measure-preserving group isomorphism \(\pi_1\) between \(B_{\Gamma}(0, r)\) and \(B_{\tilde{\Gamma}|R_q|(0, r)}\). The latter integral can readily be verified to be finite; see for instance, Proposition 5.2 of [18].

For the second item, we note that by (2.2), we have
\[
\|f\|_{X^s}^2 = \int_{B_{\Gamma}(0, r)} \mu^2(\xi) |\hat{f}(\xi)|^2 \, d\xi + \int_{\Gamma \setminus B_{\Gamma}(0, r)} (\xi)^2 |\hat{f}(\xi)|^2 \, d\xi. \tag{2.22}
\]

Then by Cauchy-Schwarz, (2.18), and (2.22) we have
\[
\int_{B_{\Gamma}(0, r)} |\hat{f}(\xi)| \, d\xi \leq \left(\int_{B_{\Gamma}(0, r)} \frac{1}{\mu^2(\xi)} \, d\xi\right)^{1/2} \left(\int_{B_{\Gamma}(0, r)} \mu^2(\xi) |\hat{f}(\xi)|^2 \, d\xi\right)^{1/2} \lesssim \|f\|_{X^s}. \tag{2.23}
\]

Then (2.19) follows immediately from (2.22) and (2.23).
\[\Box\]

The next theorem records the fundamental completeness and embedding properties of the space \(X^s(\Gamma; \mathbb{F})\) in the first and second cases considered in Lemma 2.2.
where

\[ \hat{f} = \mathcal{F}^{-1}[\hat{f}] \in C_0^\infty(\Sigma; \mathbb{C}) \quad \text{and} \quad \hat{f}_{l,t} = \mathcal{F}^{-1}[\mathbf{1}_{B(0,t)} \mathcal{F}[f]] \in H^s(\Sigma; \mathbb{C}). \]

Furthermore, we have the estimates

\[ \|f_{l,t}\|_{C_0^k} = \sum_{|\alpha| \leq k} \|\partial^\alpha f_{l,t}\|_{L^\infty} \lesssim \|f_{l,t}\|_{X^s} \quad \text{and} \quad \|f_{h,t}\|_{H^s} \lesssim \|f_{h,t}\|_{X^s}. \]

(2) (Supercritical specialized Sobolev times standard Sobolev) If \( s > d/2 \), then for any \( f \in X^s(\Gamma; \mathbb{F}) \), \( g \in H^s(\Sigma; \mathbb{F}) \) we have \( fg \in H^s(\Sigma; \mathbb{F}) \), and there exists a constant \( c = c(d,s) > 0 \) for which

\[ \|fg\|_{H^s} \leq c \|f\|_{X^s} \|g\|_{H^s} \quad \text{for all} \quad f \in X^s(\Sigma; \mathbb{F}) \quad \text{and} \quad g \in H^s(\Sigma; \mathbb{F}). \]

(3) (\( \partial_1 \)-derivatives of specialized Sobolev are \( \hat{H}^{-1} \) bounded) If \( s \geq 1 \), then there exists a constant \( c = c(d,s) > 0 \) such that

\[ \|\partial_1 \eta\|_{H^{-1}} \leq c \|\eta\|_{X^s} \quad \text{for all} \quad f \in X^s(\Sigma; \mathbb{F}). \]

This implies that the map \( \partial_1 : X^s(\Sigma; \mathbb{F}) \to \hat{H}^{-1}(\Sigma; \mathbb{F}) \cap H^{-1}(\Sigma; \mathbb{F}) \) is continuous and injective.

Proof. If \( R_\Sigma = \emptyset \) or \( R_\Sigma = \{1\} \), then by the first item of Lemma 2.2, \( X^s(\Sigma; \mathbb{R}) \) is the standard Sobolev space \( H^s(\Sigma; \mathbb{R}) \). All three items then follow from standard results on \( H^s(\Sigma; \mathbb{R}) \). Suppose \( 1 \in R_\Sigma \) and \( |R_\Sigma| \geq 2 \). We note that by Lemma A.6 and Remark 3.2, the calculations performed over \( \mathbb{R}^d \) in [18] are also valid over the low frequencies belonging to general \( \hat{\Gamma} \), therefore all three items follow from minimal modifications of Theorem 5.5, the last item of Theorem 5.6, and Theorem 5.12 of [18].

2.2. The anisotropic space \( X^s(\Sigma; \mathbb{F}) \) as an algebra. Let \( \Sigma \) be defined as in (1.1). In this subsection we prove that \( X^s(\Sigma; \mathbb{F}) \) is an algebra for \( s > \frac{d}{2} \). To prove this we will first adapt the anisotropic Littlewood-Paley techniques used in [11] to prove that \( X^s(\mathbb{R}^d; \mathbb{C}) \) is an algebra. This special case turns out to be sufficient for deducing the result in the general case when \( \mathbb{R}^d \) is replaced by \( \Sigma \) and \( \mathbb{C} \) is replaced by \( \mathbb{R} \).

First, recall that the multiplier \( \mu \) defined in (2.4) satisfies \( 1/\mu \in L^2(B(0,1)) \) by the first item of Lemma 2.4. Next, we consider the functional

\[ I[F,G,H] = \int_{B(0,1)^2} \frac{\mu(\xi + \eta)}{\mu(\xi)\mu(\eta)} F(\xi)G(\eta)H(\xi + \eta) \, d\xi \, d\eta, \]

(2.28)
where $L^0(\mathbb{R}^d; [0, \infty))$ denote the non-negative measurable functions on $\mathbb{R}^d$. Our goal is to use the same formula (2.28) to define a trilinear functional over $(L^2(\mathbb{R}^d; \mathbb{C}))^3$, but for now we only define $I$ over non-negative measurable functions so that $I$ is clearly well-defined.

The next lemma shows that $I$ induces a bounded trilinear functional over $(L^2(\mathbb{R}^d; \mathbb{C}))^3$ into $\mathbb{C}$ as long as $I$ is bounded over $(L^0(\mathbb{R}^d; [0, \infty]))^3$ into $\mathbb{R}$.

**Lemma 2.7.** Suppose there exists a constant $c > 0$ such that

$$I[F, G, H] \leq c \|F\|_{L^2} \|G\|_{L^2} \|H\|_{L^2}$$

where $I$ is the functional defined in (2.28). Then $I$ induces a bounded trilinear map defined over $(L^2(\mathbb{R}^d; \mathbb{C}))^3$ into $\mathbb{C}$ via the same formula in (2.28). Furthermore, there exists a constant $C > 0$ such that

$$|I[F, G, H]| \leq C \|F\|_{L^2} \|G\|_{L^2} \|H\|_{L^2}$$

for all $F, G, H \in L^2(\mathbb{R}^d; \mathbb{C})$.

**Proof.** We first consider the case when $F, G, H \in (L^2(\mathbb{R}^d; \mathbb{R}))^3$. Then by decomposing $F = F_0 - F_1$ where $F_0 = F^+$ and $F_1 = F^-$ are the positive and negative parts of $F$, and similarly for $G, H$, we have

$$FGH = \sum_{j,k,l \in \{0,1\}} (-1)^{j+k+l}F_jG_kH_l.$$  \hfill (2.31)

By linearity,

$$I[F, G, H] = \sum_{j,k,l \in \{0,1\}} (-1)^{j+k+l}I[F_j, G_k, H_l].$$  \hfill (2.32)

By assumption, $|I[F_j, G_k, H_l]| \leq c \|F_j\|_{L^2} \|G_k\|_{L^2} \|H_l\|_{L^2}$ for all $j, k, l \in \{0,1\}$, thus by (2.32) $I$ is bounded over $F, G, H \in (L^2(\mathbb{R}^d; \mathbb{R}))^3$ and $|I[F, G, H]| \leq C \|F\|_{L^2} \|G\|_{L^2} \|H\|_{L^2}$ for all $F, G, H \in L^2(\mathbb{R}^d; \mathbb{R})$.

In the general case when $F, G, H \in (L^2(\mathbb{R}^d; \mathbb{C}))^3$, we write $F = F_0 + iF_1$ where $F_0, F_1 \in L^2(\mathbb{R}^d; \mathbb{R})$, and similarly for $G, H$. Then

$$FGH = \sum_{j,k,l \in \{0,1\}} (i)^{j+k+l}F_jG_kH_l,$$  \hfill (2.33)

and by linearity again we have

$$I[F, G, H] = \sum_{j,k,l \in \{0,1\}} (i)^{j+k+l}I[F_j, G_k, H_l].$$  \hfill (2.34)

By the first case and following the same line of reasoning, $I$ is bounded over $(L^2(\mathbb{R}^d; \mathbb{C}))^3$ into $\mathbb{C}$ and there exists a constant $C > 0$ for which (2.30) holds.

Next, we claim that supercritical specialized Sobolev space is an algebra if and only if the functional defined via (2.28) is bounded over $(L^0(\mathbb{R}^d; [0, \infty]))^3$ into $[0, \infty]$.

**Proposition 2.8.** Assume $s > d/2$. Then $X^s(\mathbb{R}^d; \mathbb{C})$ is an algebra if and only if the mapping $I$ defined via (2.28), there exists a constant $c > 0$ such that for all $F, G, H \in L^2(\mathbb{R}^d; [0, \infty])$ we have the estimate $|I[F, G, H]| \leq c \|F\|_{L^2} \|G\|_{L^2} \|H\|_{L^2}$.

**Proof.** Assume first that $I$ is bounded over $(L^2(\mathbb{R}^d; [0, \infty]))^3$. By Lemma 2.7, $I$ induces a bounded trilinear functional over $(L^2(\mathbb{R}^d; \mathbb{C}))^3$ into $\mathbb{C}$ via the same formula (2.28), and there exists a constant $C > 0$ such that

$$|I[F, G, H]| \leq C \|F\|_{L^2} \|G\|_{L^2} \|H\|_{L^2}.$$  \hfill (2.35)

for all $F, G, H \in L^2(\mathbb{R}^d; \mathbb{C})$.

Now suppose that $f, g \in X^s(\mathbb{R}^d; \mathbb{C})$. By the second item in Theorem 2.6, we may write $f = f_0 + f_1$ where $f_0 = \mathcal{F}^{-1}[1_{B(0,1)}\mathcal{F}[f]] \in C_0^\infty(\mathbb{R}^d; \mathbb{C})$ and $f_1 = \mathcal{F}^{-1}[1_{\mathbb{R}^d \setminus B(0,1)}\mathcal{F}[f]] \in H^s(\mathbb{R}^d; \mathbb{C})$, and similarly for $g$. Then we have the decomposition

$$fg = f_0g_0 + f_0g_1 + f_1g_0 + f_1g_1.$$  \hfill (2.36)

If $1 \leq i + j$, by the fact that we are in the supercritical regime $s > d/2$, $f_1, g_1 \in H^s(\mathbb{R}^d; \mathbb{C})$, (2.25), and (2.26) we have $f_1g_1 \in H^s$ with the estimate

$$\|f_1g_1\|_{H^s} \lesssim \|f_0\|_{X^s} \|g_0\|_{X^s}.$$  \hfill (2.37)

Thus, it remains to use the boundedness of $I$ to understand the product $f_0g_0$. Note that by the first item of Theorem 2.6, we have the inclusions $f_0, g_0 \in L^1(\mathbb{R}^d; \mathbb{C})$, and hence Young’s inequality implies that $f_0 * g_0 \in L^1(\mathbb{R}^d; \mathbb{C})$; also, $\text{supp}(f_0 * g_0) \subseteq B(0, 2)$. 


Now let \( \varphi \in L^2(\mathbb{R}^d; \mathbb{C}) \cap L^\infty(\mathbb{R}^d; \mathbb{C}) \). Since \( \hat{f}_0 \) and \( \hat{g}_0 \) are supported in \( B(0, 1) \) and \( \mu \) is locally bounded, we may employ Tonelli’s theorem and a change of variables to see that

\[
\int_{\mathbb{R}^d} \mu(\hat{f}_0 * \hat{g}_0) \varphi = \int_{\mathbb{R}^d} \int_{\mathbb{R}^d} \mu(\xi + \eta) \hat{f}_0(\xi) \hat{g}_0(\eta) \varphi(\xi + \eta) \, d\xi \, d\eta = \int_{B(0,1)^2} \mu(\xi + \eta) \hat{f}_0(\xi) \hat{g}_0(\eta) \varphi(\xi + \eta) \, d\xi \, d\eta = I[\mu \hat{f}_0, \mu \hat{g}_0, \varphi].
\]  

(2.38)

Therefore, by (2.35) we have

\[
\left| \int_{\mathbb{R}^d} \mu(\hat{f}_0 * \hat{g}_0) \varphi \right| \leq C \| \mu \hat{f}_0 \|_{L^2} \| \mu \hat{g}_0 \|_{L^2} \| \varphi \|_{L^2},
\]  

(2.39)

but by the density of \( L^2(\mathbb{R}^d; \mathbb{C}) \cap L^\infty(\mathbb{R}^d; \mathbb{C}) \) in \( L^2(\mathbb{R}^d; \mathbb{C}) \), the left side of this expression extends to define a bounded linear functional on \( L^2(\mathbb{R}^d; \mathbb{C}) \) obeying the same estimate. Upon invoking Riesz’s representation theorem, we deduce that \( \mu(\hat{f}_0 * \hat{g}_0) \in L^2(\mathbb{R}^d; \mathbb{C}) \) with

\[
\| \mu(\hat{f}_0 * \hat{g}_0) \|_{L^2} \leq C \| \mu \hat{f}_0 \|_{L^2} \| \mu \hat{g}_0 \|_{L^2} \leq C \| f \|_{X^*} \| g \|_{X^*}.
\]  

(2.40)

As \( \text{supp}(\hat{f}_0 * \hat{g}_0) \subseteq B(0, 2) \), we have that \( \left\| \langle \cdot \rangle^{-1} \mu(\hat{f}_0 * \hat{g}_0) \right\|_{L^2} \leq \frac{5}{2^d} \left\| \mu(\hat{f}_0 * \hat{g}_0) \right\|_{L^2}. \) Thus, using (2.37), (2.40), and the fourth item of Theorem 2.6 we find that

\[
\| fg \|_{X^*} \leq \sum_{0 \leq i+j \leq 2} \| f_i g_j \|_{X^*} = \| \langle \cdot \rangle^{i+j} (\hat{f}_0 * \hat{g}_0) \|_{L^2} + \sum_{1 \leq i+j \leq 2} \| f_i g_j \|_{X^*} \leq \| \mu(\hat{f}_0 * \hat{g}_0) \|_{L^2} + \sum_{1 \leq i+j \leq 2} \| f_i g_j \|_{H^*} \leq \| f \|_{X^*} \| g \|_{X^*}. \]  

(2.41)

Thus, \( X^* (\mathbb{R}^d; \mathbb{C}) \) is an algebra.

Conversely, assume that \( X^* (\mathbb{R}^d; \mathbb{C}) \) is an algebra. Let \( F, G, H \in L^2(\mathbb{R}^d; [0, \infty)) \) and observe first that \( I[F, G, H] = I[F^1 B(0,1), G^1 B(0,1), H] \) and then that \( \mathcal{F}^{-1} [F^1 B(0,1)/|\mu|, \mathcal{F}^{-1} [G^1 B(0,1)/|\mu|] \in X^* (\mathbb{R}^d; \mathbb{C}) \). Therefore, by Cauchy-Schwarz and the boundedness of products in \( X^* \) we have

\[
I[F, G, H] = I[F^1 B(0,1), G^1 B(0,1), H] = \int_{\mathbb{R}^d} \mu((F^1 B(0,1)/|\mu|) * (G^1 B(0,1)/|\mu|) H)
\leq \| \mu((F^1 B(0,1)/|\mu|) * (G^1 B(0,1)/|\mu|) H) \|_{L^2} \| H \|_{L^2} = \| \mathcal{F}^{-1} [F^1 B(0,1)/|\mu|, \mathcal{F}^{-1} [G^1 B(0,1)/|\mu|] X^* \| \| H \|_{L^2} \leq \| F \|_{L^2} \| G \|_{L^2} \| H \|_{L^2}.
\]  

(2.42)

Hence, \( I \) is bounded over \( (L^2(\mathbb{R}^d; [0, \infty)))^3 \) and the proof is complete. \( \square \)

Thus, by Lemma 2.7 and Proposition 2.8, to prove that \( X^* (\mathbb{R}^d; \mathbb{C}) \) is an algebra it remains to show that the functional \( I \) defined over \( (L^2(\mathbb{R}^d; [0, \infty]))^3 \) via (2.28) satisfies (2.29). For the rest of this subsection we assume that \( F, G, H \in L^2(\mathbb{R}^d; [0, \infty)) \), and we now estimate the \( L^2 \)-boundedness of the operator \( I \).

First, we introduce a decomposition of the frequency space \( \mathbb{R}^d \times \mathbb{R}^d \).

**Definition 2.9** (Squared frequency space and functional decomposition). We write \( I \) as the sum of two operators \( I = I_0 + I_1 \), where each \( I_i \) is accounting for the contribution from a special portion of squared frequency space.

1. Identify the following ‘good’ and ‘bad’ sets. First we define the ‘good’ set \( E_0 \) via

\[
E_0 = \{ (\xi, \eta) \in B(0, 1)^2 : |\xi| + |\eta| \leq 3|\xi| - |\eta| \}.
\]  

(2.43)

Next we define the ‘bad’ set \( E_1 \) via

\[
E_1 = \{ (\xi, \eta) \in B(0, 1)^2 : |\xi| + |\eta| > 3|\xi| - |\eta| \}.
\]  

(2.44)

Note that \( B(0,1)^2 = E_0 \cup E_1 \).

2. For \( i \in \{0, 1\} \), we define the functionals \( I_i : (L^2(\mathbb{R}^d; [0, \infty]))^3 \rightarrow [0, \infty] \) via

\[
I_i[F, G, H] = \int_{E_i} \frac{\mu(\xi + \eta)}{\mu(\xi)} F(\xi) G(\eta) H(\xi + \eta) \, d\xi \, d\eta.
\]  

(2.45)

Clearly, \( I_i \) is well-defined, and we have the identity \( I = I_0 + I_1 \).

Next, we analyze the set \( E_1 \) as defined in (2.44).

**Lemma 2.10.** The inclusion \( (\xi, \eta) \in E_1 \) is equivalent to the bounds \( \frac{1}{2} |\eta| < |\xi| < 2|\eta| \), and for \( (\xi, \eta) \in E_1 \) we have the estimate \( |\xi + \eta| < 3|\xi| \).
Proof. The second bound follows from the equivalence and the triangle inequality, so it suffices to prove the equivalence. For this we note that
\[
(\xi, \eta) \in E_1 \iff |\xi| + |\eta| > 3|\xi| - |\eta| \iff |\xi|^2 - \frac{5}{2}|\xi||\eta| + |\eta|^2 < 0 \iff \frac{1}{2}|\eta| < |\xi| < 2|\eta|.
\]
(2.46)

By Lemma 2.10, one can get a good geometric understanding of \(E_0\) and \(E_1\) by examining the plot in \((|\xi|, |\eta|)\) space given in Figure 3.

![Figure 3. A schematic breakdown of the sets \(E_0\) and \(E_1\)](image)

The region \(E_0\) is dubbed the good region because of the following lemma.

**Lemma 2.11** (Multiplier subadditivity in \(E_0\)). If \((\xi, \eta) \in E_0\), then \(\mu(\xi + \eta) \leq 3(\mu(\xi) + \mu(\eta))\).

**Proof.** If \(|\xi| + |\eta| \leq 3||\xi| - |\eta||\), then the reverse triangle inequality \(|\xi + \eta| \geq ||\xi| - |\eta||\) allows us to deduce
\[
\mu(\xi + \eta) = \frac{|\xi_1 + \eta_1|}{|\xi + \eta|} + |\xi + \eta| \leq \frac{|\xi_1| + |\eta_1|}{||\xi| - |\eta||} + |\xi| + |\eta| \leq 3\frac{|\xi_1| + |\eta_1|}{|\xi| + |\eta|} + |\xi| + |\eta| \leq 3(\mu(\xi) + \mu(\eta)).
\]
(2.47)

This gives us the boundedness of \(I_0\).

**Proposition 2.12.** For all \(F, G, H \in L^2(\mathbb{R}^d; [0, \infty))\) we have the estimate
\[
I_0[F, G, H] \leq 6 \|1/\mu\|_{L^2(B(0, 1))} \|F\|_{L^2} \|G\|_{L^2} \|H\|_{L^2}.
\]
(2.48)

**Proof.** By applying Lemma 2.11, Tonelli’s theorem, and Cauchy-Schwarz, we find
\[
I_0[F, G, H] \leq 3 \int_{E_0} \frac{1}{\mu(\xi)} F(\xi) G(\xi) H(\xi + \eta) \, d\xi \, d\eta + 3 \int_{E_0} \frac{1}{\mu(\eta)} F(\xi) G(\eta) H(\xi + \eta) \, d\xi \, d\eta
\]
\[
\leq 3 \int_{B(0, 1)^2} \frac{1}{\mu(\xi)} F(\xi) G(\xi) H(\xi + \eta) \, d\xi \, d\eta + 3 \int_{B(0, 1)^2} \frac{1}{\mu(\eta)} F(\xi) G(\eta) H(\xi + \eta) \, d\xi \, d\eta
\]
\[
\leq 3 \|H\|_{L^2} \left( \|G\|_{L^2} \int_{B(0, 1)} \frac{1}{\mu(\xi)} F(\xi) \, d\xi + \|F\|_{L^2} \int_{B(0, 1)} \frac{1}{\mu(\eta)} G(\eta) \, d\eta \right)
\]
\[
\leq 6 \|1/\mu\|_{L^2(B(0, 1))} \|F\|_{L^2} \|G\|_{L^2} \|H\|_{L^2}.
\]
(2.49)

We now further decompose the set \(E_1\) and the functional \(I_1\), as defined in (2.44) and (2.45).

**Definition 2.13.** We make the following definitions for \(m, n \in \mathbb{N}\).

1. We define the set
\[
E_{m,n} = \{(\xi, \eta) \in E_1 : 2^{-m-1} \leq |\xi| \leq 2^{-m} \quad \text{and} \quad 2^{-n+1} \leq |\xi + \eta| \leq 2^{-n+2}\}.
\]
(2.50)

2. We define the functional \(I_{m,n} : (L^0(\mathbb{R}^d; \mathbb{R}))^3 \to [0, \infty]\) via
\[
I_{m,n}[F, G, H] = \int_{E_{m,n}} \frac{\mu(\xi + \eta)}{\mu(\xi)\mu(\eta)} F(\xi) G(\eta) H(\xi + \eta) \, d\xi \, d\eta.
\]
(2.51)
The dyadic decomposition of $E_1$ gives us the following.

**Lemma 2.14.** The following hold.

1. We have the equalities
   \[ \sum_{m=0}^{\infty} \sum_{n=m}^{\infty} E_{m,n} = E_1 \quad \text{and} \quad \sum_{m=0}^{\infty} \sum_{n=m}^{\infty} I_{m,n} = I_1. \]  
   \[ (2.52) \]

2. Let $F, G, H \in L^2(\mathbb{R}^d, [0, \infty])$. Then
   \[ I_{m,n}[F, G, H] \leq I_1[F_{m}, G_{m}, H_n], \]
   for the functions
   \[ F_m = F\mathbb{1}_{2^{-m-1}A}, \quad G_m = G\mathbb{1}_{2^{-m-1}A}, \quad \text{and} \quad H_n = H\mathbb{1}_{2^{-n}A}, \]
   where $A$ is the annulus $A = \{ x \in \mathbb{R}^d : 1/2 \leq |x| \leq 4 \}$.
   \[ (2.53) \]

3. We have $L^d$-almost everywhere the inequality
   \[ \sum_{m \in \mathbb{N}} \mathbb{1}_{2^{-m-1}A} \leq 5\mathbb{1}_{B(0,2)} \]
   \[ (2.55) \]

**Proof.** For the first item we note that the second estimate of Lemma 2.10 implies that if $n < m$ then $E_{m,n} = \emptyset$, and hence $I_{m,n} = 0$. The second item follows from the first item of Lemma 2.10 and nonnegativity. To prove the third item, we note that if $x \in B(0, 2) \setminus \{0\}$, there exists a unique $n \in \mathbb{N}$ such that $|x| \in 2^{-n}[1, 2)$. It is clear that $\mathbb{1}_{2^{-m-1}A}(x) = 0$ whenever $m \geq 0$ and $|m - n| > 2$. The inequality now follows.

Using the preceding lemma we arrive at the following proposition.

**Proposition 2.15.** Let $d \in \mathbb{N}$. The following hold:

1. If $d \geq 4$, there exists a constant $C_d > 0$ such that for all $F, G, H \in L^2(\mathbb{R}^d, [0, \infty])$ we have the estimate
   \[ I_1[F, G, H] \leq C_d \| F \|_{L^2} \| G \|_{L^2} \| H \|_{L^2}. \]
   \[ (2.56) \]

2. If $d \in \{2, 3\}$, there exists a constant $C_{1,d} > 0$ such that for all $F, G, H \in L^2(\mathbb{R}^d, [0, \infty])$ we have the estimate
   \[ \sum_{m=0}^{\infty} \sum_{n=2m+1}^{\infty} I_{m,n}[F, G, H] \leq C_{1,d} \| F \|_{L^2} \| G \|_{L^2} \| H \|_{L^2}. \]
   \[ (2.57) \]

**Proof.** Let $m, n, d \in \mathbb{N}$ with $2 \leq d$ and $m \leq n$. According to the second item of Lemma 2.14, we have

\[ I_{m,n}[F, G, H] \leq I_1[F_{m}, G_{m}, H_n] = \int_{E_{1}} \frac{\mu(\xi + \eta)}{\mu(\xi) \mu(\eta)} F_m(\xi) G_m(\eta) H_n(\xi + \eta) \, d\xi \, d\eta. \]

For $(\xi, \eta) \in E_1$, the right hand integrand vanishes except possibly when the following inequalities hold: $2^{-m-2} \leq |\xi| \leq 2^{-m+1}$, $2^{-m-2} \leq |\eta| \leq 2^{-m+1}$, and $2^{-n-1} \leq |\xi + \eta| \leq 2^{-n+2}$. This provides the elementary estimate

\[ \frac{\mu(\xi + \eta)}{\mu(\xi) \mu(\eta)} \leq \frac{1 + |\xi + \eta|}{|\xi| |\eta|} \leq 2^{2m+4} (1 + 2^{n+2}) \lesssim 2^{2m}, \]

and hence

\[ I_{m,n}[F, G, H] \lesssim 2^{2m} \int_{E_1} F_m(\xi) G_m(\eta) H_n(\xi + \eta) \, d\xi \, d\eta. \]

Now, for $\ell \in \mathbb{Z}^d$ let $Q_\ell$ denote the closed cube centered at $2^{-n}\ell$ of side length $2^{-n}$ and $\tilde{Q}_\ell$ denote the closed cube centered at $-2^{-n}\ell$ of side length $9 \cdot 2^{-n}$. By construction, almost everywhere we have that

\[ 1 = \sum_{\ell \in \mathbb{Z}^d} \mathbb{1}_{Q_\ell} \quad \text{and} \quad 9^d = \sum_{\ell \in \mathbb{Z}^d} \mathbb{1}_{\tilde{Q}_\ell}, \]

and we may compute

\[ \max_{\ell \in \mathbb{Z}^d} \| \mathbb{1}_{Q_\ell} \|_{L^2} = \| \mathbb{1}_{Q_0} \|_{L^2} = 2^{-nd/2}. \]

Notice that if $(\xi, \eta) \in E_1$ is such that the integrand on the right hand side of (2.59) is nonzero and $\xi \in Q_\ell$, then

\[ |\eta + 2^{-n}\ell|_\infty \leq |\eta + \xi| + | - \xi + 2^{-n}\ell|_\infty \leq 2^{-n+2} + 2^{-n-1} = \frac{9}{2} \cdot 2^{-n}, \]

which in particular implies that $\eta \in \tilde{Q}_\ell$. Hence,

\[ \int_{E_1} F_m(\xi) G_m(\eta) H_n(\xi + \eta) \, d\xi \, d\eta \lesssim \sum_{\ell \in \mathbb{Z}^d} \int_{B(0,1)^2} F_m(\xi) \mathbb{1}_{Q_\ell}(\xi) G_m(\eta) \mathbb{1}_{\tilde{Q}_\ell}(\eta) H_n(\xi + \eta) \, d\xi \, d\eta. \]

\[ (2.63) \]
By Tonelli’s theorem, repeated applications of Cauchy-Schwarz, and the identities in (2.60) and (2.61) we then find that
\[\sum_{\ell \in \mathbb{Z}^d} \int_{B(0,1)^2} F_m(\xi) \mathbb{1}_{Q_\ell}(\xi) G_m(\eta) \mathbb{1}_{\tilde{Q}_\ell}(\eta) H_n(\xi + \eta) \, d\xi \, d\eta\]
\[= \sum_{\ell \in \mathbb{Z}^d} \int_{B(0,1)} F_m(\xi) \mathbb{1}_{Q_\ell}(\xi) \int_{B(0,1)} G_m(\eta) \mathbb{1}_{\tilde{Q}_\ell}(\eta) H_n(\xi + \eta) \, d\eta \, d\xi\]
\[\leq \sum_{\ell \in \mathbb{Z}^d} \int_{B(0,1)} F_m(\xi) \mathbb{1}_{Q_\ell}(\xi) \|G_m\|_{L^2} \|H_n\|_{L^2} \, d\xi \leq \|H_n\|_{L^2} \sum_{\ell \in \mathbb{Z}^d} \|F_m\|_{L^2} \|G_m\|_{L^2} \|H_n\|_{L^2}\]
\[\leq 2^{-nd/2} \|H_n\|_{L^2} \left( \int_{\mathbb{R}^d} |F_m(\xi)|^2 \, d\xi \right)^{1/2} \left( \int_{\mathbb{R}^d} |G_m(\eta)|^2 \, d\eta \right)^{1/2} \|H_n\|_{L^2}\]
\[\lesssim 2^{-nd/2} \|F_m\|_{L^2} \|G_m\|_{L^2} \|H_n\|_{L^2}. \quad (2.64)\]

Next we combine equations (2.59), (2.63), and (2.64) to deduce that
\[I_{m,n}[F,G,H] \lesssim 2^{m-n/2} \|F\|_{L^2} \|G\|_{L^2} \|H\|_{L^2}. \quad (2.65)\]

Now we break to cases based on the dimension. First, if \(d \geq 4\), then there is no issue in summing over the set \(\{(m,n) \in \mathbb{N}^2 : m \leq n\}\). Indeed, we use Cauchy-Schwarz along with the third item of Lemma 2.14 to find
\[\sum_{m=0}^\infty \sum_{n=m}^\infty I_{m,n}[F,G,H] \lesssim \sum_{m=0}^\infty \sum_{n=m}^\infty 2^{m-n/2} \|F_m\|_{L^2} \|G_m\|_{L^2} \|H_n\|_{L^2}\]
\[\leq \sum_{m=0}^\infty 2^{2m} \|F_m\|_{L^2} \|G_m\|_{L^2} \left( \sum_{n=m}^\infty 2^{-n} \right)^{1/2} \left( \int_{\mathbb{R}^d} |H(\xi)|^2 \, d\xi \right)^{1/2}\]
\[\lesssim \|H\|_{L^2} \sum_{m=0}^\infty 2^{(2-d/2)m} \|F_m\|_{L^2} \|G_m\|_{L^2} \lesssim \|F\|_{L^2} \|G\|_{L^2} \|H\|_{L^2}. \quad (2.66)\]

This completes the proof of the first item.

For the second item, we let \(d \in \{2,3\}\) in equation (2.65) and sum over the set \(\{(m,n) \in \mathbb{N}^2 : 2m+1 \leq n\}\). Arguing as above, we arrive at the bound
\[\sum_{m=0}^\infty \sum_{n=2m+1}^\infty I_{m,n}[F,G,H] \lesssim \|H\|_{L^2} \sum_{m=0}^\infty 2^{2m} \left( \sum_{n=2m+1}^\infty 2^{-d(n-2m)} \right)^{1/2} \|F_m\|_{L^2} \|G_m\|_{L^2}\]
\[\lesssim \|F\|_{L^2} \|G\|_{L^2} \|H\|_{L^2}. \quad (2.67)\]

This completes the proof of the second item. \(\square\)

It remains to study the functional \(\sum_{m=0}^\infty \sum_{n=2m}^\infty I_{m,n}\) in the case \(d \in \{2,3\}\). We do this in the subsequent proposition.

**Proposition 2.16.** Let \(d \in \{2,3\}\). There exists a constant \(C_{2,d} > 0\) such that for all \(F,G,H \in L^2(\mathbb{R}^d;[0,\infty))\) we may estimate
\[\sum_{m=0}^\infty \sum_{n=2m}^\infty I_{m,n}[F,G,H] \leq C_{2,d} \|F\|_{L^2} \|G\|_{L^2} \|H\|_{L^2}. \quad (2.68)\]

**Proof.** Let \(m,n \in \mathbb{N}\) with \(m \leq n \leq 2m\). We begin by estimating the right hand side of the following inequality of Lemma 2.14, \(I_{m,n}[F,G,H] \leq I_1[F_m,G_m,H_n]\), by decomposing into a rectangular grid. To this end we define the following family of rectangles indexed by \(p \in \mathbb{Z}, \pi \in \mathbb{Z}^{d-1}, \) and \(\alpha \in [1,\infty), \)
\[R_{p,\pi}(\alpha) = 2^{-2m}[-\alpha/2 + p, \alpha/2 + p] \times \prod_{\nu=1}^{d-1} [-\alpha/2 + \pi_\nu, \alpha/2 + \pi_\nu]. \quad (2.69)\]
Suppose that \((\xi, \eta) \in E_1\) are such that the integrand of \(I_1[F_m,G_m,H_n]\) is nonzero and that additionally \(\xi \in R_{p,\pi}(1)\) and \(\eta \in R_{q,\sigma}(1)\) for some \(p,q \in \mathbb{Z}\) and \(\pi,\sigma \in \mathbb{Z}^{d-1}\). Then we obtain the inequalities
\[2^{-n}(\pi_\nu + \sigma_\nu - 1/2) \leq \xi_{\nu+1} + \eta_{\nu+1} \leq 2^{-n}(\pi_\nu + \sigma_\nu + 1/2) \text{ for all } 1 \leq \nu \leq d-1, \quad (2.70)\]
and upon pairing these with the bound $|\xi + \eta| \leq 2^{-n+2}$ we deduce that
\[
|\pi_\nu + \sigma_\nu| \leq |\pi_\nu + \sigma_\nu - 2^n(\xi_{\nu+1} + \eta_{\nu+1})| + 2^n|\xi_{\nu+1} + \eta_{\nu+1}| \leq 9/2 \text{ for all } 1 \leq \nu \leq d - 1. \tag{2.71}
\]

In particular, this implies that $\eta \in R_{q,-\pi}(9)$, which when combined with the inclusion $\xi \in R_{p,\pi}(1)$ further implies that $\xi + \eta \in R_{p+q,0}(10)$. Therefore, we may estimate
\[
I_{m,n}[F, G, H] \leq \sum_{p,q \in \mathbb{Z}^d} \sum_{\pi \in \mathbb{Z}^{d-1}} \int_{E_1} \frac{\mu(\xi + \eta)}{\mu(\xi \mu(\eta))} F_m(\xi) \mathbbm{1}_{R_p,\pi}(1)(\xi) G_m(\eta) \mathbbm{1}_{R_q,\pi}(9)(\eta) H_n(\xi + \eta) \mathbbm{1}_{R_{p+q,0}(10)}(\xi + \eta) \, d\xi \, d\eta.
\]

For fixed $p, q,$ and $\pi$, we now want to estimate the factor $\mu(\xi + \eta)/\mu(\xi)\mu(\eta)$ appearing in the integrand in (2.72). First we note that if $\xi \in R_{p,\pi}(1) \cap 2^{-m-1}A$, then $|2^m\xi - p| \leq 1/2$, which implies that $2^m|\xi| \geq \max \{0, |p| - 1/2\}$ and hence
\[
\mu(\xi) = \frac{|\xi|}{|\xi|} + |\xi| \geq 2^{-m-1}(2^m|\xi| + 1/2) \geq 2^{-m} \max \{1, |p|\}. \tag{2.73}
\]

Similarly, for $\eta \in R_{q,-\pi}(9) \cap 2^{-m-1}A$, we have the bound
\[
\mu(\eta) \geq 2^{-m-1}(\max \{|q| - 9/2, 0\} + 1/2) \geq 2^{-m} \max \{1, |q|\}. \tag{2.74}
\]

Finally, for $\xi + \eta \in R_{p+q,0}(10) \cap 2^{-n}A$, we estimate $\mu(\xi + \eta)$ from above, recalling that $m \leq n$:
\[
\mu(\xi + \eta) \leq 2^{n+1}|\xi_1 + \eta_1| + 2^{-n+2} = 2^{n-2m+1}(2^m|\xi_1 + \eta_1| + 2^{(m-n)+1}) \\
\leq 2^{n-2m+1}(|p + q| + 5/2) \lesssim 2^{n-2m}(\max \{1, |p|\} + \max \{1, |q|\}). \tag{2.75}
\]

Upon piecing equations (2.73), (2.74), and (2.75) all together, we obtain the bound
\[
\frac{\mu(\xi + \eta)}{\mu(\xi)\mu(\eta)} \lesssim 2^n \left(\frac{1}{\max \{1, |p|\}} + \frac{1}{\max \{1, |q|\}}\right) \tag{2.76}
\]
for $\xi \in R_{p,\pi}(1) \cap 2^{-m-1}A$ and $\eta \in R_{q,-\pi}(9) \cap 2^{-m-1}A$.

Next, we combine (2.72) and (2.76) and then use Tonelli’s theorem and Cauchy-Schwarz to obtain
\[
2^{-n}I_{m,n}[F, G, H] \leq \sum_{p,q \in \mathbb{Z}^d} \sum_{\pi \in \mathbb{Z}^{d-1}} \frac{1}{\max \{1, |p|\}} \int_{B(0,1)} F_m(\xi) \mathbbm{1}_{R_p,\pi}(1)(\xi) \int_{B(0,1)} G_m(\eta) \mathbbm{1}_{R_q,\pi}(9)(\eta) H_n(\xi + \eta) \mathbbm{1}_{R_{p+q,0}(10)}(\xi + \eta) \, d\eta \, d\xi \\
+ \sum_{p,q \in \mathbb{Z}^d} \sum_{\pi \in \mathbb{Z}^{d-1}} \frac{1}{\max \{1, |q|\}} \int_{B(0,1)} F_m(\xi) \mathbbm{1}_{R_p,\pi}(1)(\xi) \int_{B(0,1)} G_m(\eta) \mathbbm{1}_{R_q,\pi}(9)(\eta) H_n(\xi + \eta) \mathbbm{1}_{R_{p+q,0}(10)}(\xi + \eta) \, d\eta \, d\xi \\
\lesssim \sum_{p,q \in \mathbb{Z}^d} \sum_{\pi \in \mathbb{Z}^{d-1}} \left(\frac{1}{\max \{1, |p|\}} + \frac{1}{\max \{1, |q|\}}\right) \int_{B(0,1)} F_m(\xi) \mathbbm{1}_{R_p,\pi}(1)(\xi) \|G_m\|_{L^2} \|H_n\|_{L^2} \|\mathbbm{1}_{R_{p+q,0}(10)}\|_{L^2}^2 \, d\xi \\
\lesssim \sum_{p,q \in \mathbb{Z}^d} \sum_{\pi \in \mathbb{Z}^{d-1}} \left(\frac{1}{\max \{1, |p|\}} + \frac{1}{\max \{1, |q|\}}\right) \|F_m\|_{L^2} \|\mathbbm{1}_{R_p,\pi}(1)\|_{L^2} \|G_m\|_{L^2} \|H_n\|_{L^2} \|\mathbbm{1}_{R_{p+q,0}(10)}\|_{L^2}^2. \tag{2.77}
\]

What remains, before we sum over $m$ and $n$, is a few more applications of Cauchy-Schwarz, one for each sum above. First we handle the sum over $\pi \in \mathbb{Z}^{d-1}$ with
\[
\sum_{\pi \in \mathbb{Z}^{d-1}} \|F_m\|_{L^2} \|G_m\|_{L^2} \|H_n\|_{L^2} \lesssim \left(\int_{\mathbb{R}^d} |F_m(\xi)|^2 \sum_{\pi \in \mathbb{Z}^{d-1}} \|\mathbbm{1}_{R_p,\pi}(1)\|_{L^2} \|G_m\|_{L^2} \|H_n\|_{L^2} \|\mathbbm{1}_{R_{p+q,0}(10)}\|_{L^2} \right)^{1/2} \tag{2.78}
\]
\[
\lesssim \|F_m\|_{L^2} \sum_{\pi \in \mathbb{Z}^{d-1}} \|\mathbbm{1}_{R_{p,\pi}(1)}\|_{L^2} \|G_m\|_{L^2} \|H_n\|_{L^2} \|\mathbbm{1}_{R_{p+q,0}(10)}\|_{L^2} \left(\int_{\mathbb{R}^d} |G_m(\eta)|^2 \sum_{\pi \in \mathbb{Z}^{d-1}} \|\mathbbm{1}_{R_q,\pi}(9)\|_{L^2} \right)^{1/2}.
\]
Next we consider the sums over \( p, q, \in \mathbb{Z} \). First, we consider the one containing \( 1/\max\{1, |p|\} \),

\[
\sum_{p,q \in \mathbb{Z}} \frac{1}{\max\{1, |p|\}} \left| I_{p,q}(1) \right| \| F_{p,q} \|_{L^2} \| G_{p,q} \|_{L^2} \| H_{p,q} \|_{L^2} \leq \left( \sum_{p \in \mathbb{Z}} 1/\max\{1, p^2\} \right)^{1/2} \left( \sum_{q \in \mathbb{Z}} 1/\max\{1, q^2\} \right)^{1/2} \| F_{p,q} \|_{L^2} \| G_{p,q} \|_{L^2} \| H_{p,q} \|_{L^2} \leq \| F \|_{L^2} \| G \|_{L^2} \| H \|_{L^2}.
\]

Similarly, we consider the one containing \( 1/\max\{1, |q|\} \),

\[
\sum_{p,q \in \mathbb{Z}} \frac{1}{\max\{1, |q|\}} \left| I_{p,q}(1) \right| \| F_{p,q} \|_{L^2} \| G_{p,q} \|_{L^2} \| H_{p,q} \|_{L^2} \leq \left( \sum_{q \in \mathbb{Z}} 1/\max\{1, q^2\} \right)^{1/2} \left( \sum_{p \in \mathbb{Z}} 1/\max\{1, p^2\} \right)^{1/2} \| F_{p,q} \|_{L^2} \| G_{p,q} \|_{L^2} \| H_{p,q} \|_{L^2} \leq \| F \|_{L^2} \| G \|_{L^2} \| H \|_{L^2}.
\]

Synthesizing equations (2.77), (2.78), (2.79), and (2.80) we find that

\[
I_{m,n} \leq 2^{-m+n(3-d)/2} \| F \|_{L^2} \| G \|_{L^2} \| H \|_{L^2}.
\]

Hence, we may sum over \( m \leq n \leq 2m \) to bound

\[
\sum_{m=0}^{2m} I_{m,n} \leq 2^{-m} \left( \sum_{m=0}^{2m} 2^{3-d} \right)^{1/2} \| F \|_{L^2} \| G \|_{L^2} \leq 2^{-m} \| F \|_{L^2} \| G \|_{L^2} \| H \|_{L^2},
\]

and finally sum over \( m \geq 0 \) to see that

\[
\sum_{m=0}^{\infty} \sum_{n=m}^{2m} I_{m,n} \leq \| F \|_{L^2} \| G \|_{L^2} \| H \|_{L^2},
\]

which is the stated estimate. \( \square \)

By synthesizing the previous results, we immediately arrive at the following result.

**Theorem 2.17 (Boundedness of I).** The functional \( I : L^2(\mathbb{R}; [0, \infty])^3 \rightarrow [0, \infty] \) defined in (2.28) is bounded and satisfies the estimate (2.30).

**Proof.** We simply combine Lemma 2.14 and Propositions 2.12, 2.15, 2.16, and 2.8. \( \square \)

With Theorem 2.17 in hand, we are ready to prove Theorem 1.1 as a corollary.

**Proof of Theorem 1.1.** By Lemma 2.7, Proposition 2.8, and Theorem 2.17, the anisotropic space \( X^s(\mathbb{R}^d; F) \) is an algebra for \( s > k/2 \) for \( F = C \) and \( F = \mathbb{R} \). Let \( \Sigma \) be defined as in (1.1) and let \( d = \dim \Sigma \). By the first item of Lemma 2.2, it suffices show that \( X^s(\Sigma; F) \) is an algebra for \( s > d/2 \) for the case when \( 1 \in R_{\Sigma} \) and \( |R_{\Sigma}| > 2 \), as the anisotropic space coincides with the standard Sobolev space otherwise. To show this we first argue similarly as in Proposition 2.8 to show that \( X^s(\Sigma; F) \) is an algebra if and only if for any \( F, G \in X^s(\Sigma; F) \), the estimate

\[
\| F_0 G_0 \|_{X^s(\Sigma)} \leq \| F_0 \|_{X^s(\Sigma)} \| G_0 \|_{X^s(\Sigma)}
\]

holds for \( F_0 = \mathcal{F}^{-1}[\mathcal{I}_{\mathcal{F}}, \mathcal{F}[F]], G_0 = \mathcal{F}^{-1}[\mathcal{I}_{\mathcal{F}}, \mathcal{F}[G]] \).

Let \( \Sigma_i, \Sigma_i^* \) be the reordered Cartesian products of \( \Sigma \) defined via (A.10), and \( \mathcal{P}_{\Sigma}, \mathcal{P}_{\Sigma}^* \) be the canonical reordering maps defined by (A.10). We first note that since \( F_0, G_0 \in X^s(\Sigma; F) \), by Lemma 2.1 the functions \( F_0, G_0 \in \mathcal{F}^{-1}[\mathcal{I}_{\mathcal{F}}, \mathcal{F}[F]], G_0 \in \mathcal{F}^{-1}[\mathcal{I}_{\mathcal{F}}, \mathcal{F}[G]] \) belong to \( X^s(\Sigma^*_i; F) \) with \( \| F_0 \|_{X^s(\Sigma^*_i)} = \| F_0 \|_{X^s(\Sigma^*_i)} \), \( \| G_0 \|_{X^s(\Sigma^*_i)} = \| G_0 \|_{X^s(\Sigma^*_i)} \), and by (A.29) we have

\[
\mathcal{F}_{\Sigma^*_i}^*[F_0, G_0] = \mathcal{F}_{\Sigma^*_i}^*[F_0, G_0].
\]
We also have \( \|F_0G_0\|_{X^r(\Sigma)} = \|\Sigma F_0G_0\|_{X^r(\Sigma)} = \|F_0,\Sigma, G_0,\Sigma\|_{X^r(\Sigma)} \), where \( \Sigma \) is defined via (A.11). So to establish (2.84) it suffices to establish the estimate

\[
\|F_0,\Sigma, G_0,\Sigma\|_{X^r(\Sigma)} \lesssim \|F_0,\Sigma\|_{X^r(\Sigma)} \|G_0,\Sigma\|_{X^r(\Sigma)}.
\]  

(2.86)

By (2.85), we have \( \text{supp} \mathcal{F}_\Sigma^+[F_0,\Sigma,], \text{supp} \mathcal{F}_\Sigma^+[G_0,\Sigma,] \subseteq B_{\Sigma}(0, r) \). Since \( \xi \in B_{\Sigma}(0, r) \iff \xi \in \hat{\Sigma} \), by Fubini’s theorem we then have

\[
\mathcal{F}_\Sigma^+[F_0,\Sigma, G_0,\Sigma,](\omega, k) = \int_{R_{\|k\|_{\Sigma}}} \sum_{m \in \Pi_{s^2} \cap \mathbb{L}} \mathcal{F}_\Sigma^+[F_0,\Sigma,](\omega - \eta, k - m) \mathcal{F}_\Sigma^+[G_0,\Sigma,](\eta, m) d\eta
\]

\[
= \int_{R_{\|k\|_{\Sigma}}} \mathcal{F}_\Sigma^+[F_0,\Sigma,](\omega - \eta, k) \mathcal{F}_\Sigma^+[G_0,\Sigma,](\eta, k) d\eta = \delta_{\eta 0} \int_{R_{\|k\|_{\Sigma}}} \mathcal{F}_\Sigma^+[F_0,\Sigma,](\omega - \eta, 0) \mathcal{F}_\Sigma^+[G_0,\Sigma,](\eta, 0) d\eta,
\]

where \( \delta_{\eta 0} \) is the Kronecker delta. This shows that \( \text{supp} \mathcal{F}_\Sigma^+[F_0,\Sigma, G_0,\Sigma,] \subseteq B_{\Sigma}(0, r) \). By Lemma A.6, \( \pi_{\Sigma} \) is an isometric measure-preserving group isomorphism between \( B_{\Sigma}(0, r) \) and \( B_{\|k\|_{\Sigma}}(0, r) \), and therefore

\[
\|F_0,\Sigma, G_0,\Sigma\|_{X^r(\Sigma)}^2 = \int_{B_{\|k\|_{\Sigma}}(0, r)} \mu^2(\xi) \langle \xi \rangle^{2(s-1)} |\mathcal{F}_\Sigma^+[F_0,\Sigma, G_0,\Sigma,](\xi)\rangle^2 d\xi
\]

\[
= \int_{B_{\|k\|_{\Sigma}}(0, r)} \mu^2(\omega) \langle \omega \rangle^{2(s-1)} |\mathcal{F}_\Sigma^+[F_0,\Sigma,](\omega, 0)\rangle^2 d\omega = \|F_0,\Sigma, (\cdot, 0)G_0,\Sigma, (\cdot, 0)\|_{X^r(R_{\|k\|_{\Sigma}})}^2.
\]

Next we note that since \( F_0,\Sigma, G_0,\Sigma \in X^s(\Sigma; F) \), we have the inclusions

\[
F_0,\Sigma,(\cdot, 0) \in \mathcal{M}(R_{\|k\|_{\Sigma}}; \mathbb{C}) \text{ and } \mathcal{F}_{R_{\|k\|_{\Sigma}}}[F_0,\Sigma,](\cdot, 0) \in L^1_{\text{loc}}(R_{\|k\|_{\Sigma}}; \mathbb{C}),
\]

and we know that \( F_0,\Sigma,(\cdot, 0) = F_0,\Sigma,(\cdot, 0) \), and \( \text{supp} \mathcal{F}_{R_{\|k\|_{\Sigma}}}[F_0,\Sigma,](\cdot, 0) \subseteq B_{\|k\|_{\Sigma}}(0, r) \). By Lemma A.6 again, we have

\[
\|F_0,\Sigma,(\cdot, 0)\|_{X^r(R_{\|k\|_{\Sigma}})}^2 = \int_{B_{\|k\|_{\Sigma}}(0, r)} \mu^2(\omega) \langle \omega \rangle^{2(s-1)} |\mathcal{F}_\Sigma^+[F_0,\Sigma,](\omega, 0)\rangle^2 d\omega
\]

\[
= \int_{B_{\|k\|_{\Sigma}}(0, r)} \mu^2(\xi) \langle \xi \rangle^{2(s-1)} |\mathcal{F}_\Sigma^+[F_0,\Sigma,](\xi)\rangle^2 d\xi = \|F_0,\Sigma\|_{X^r(\Sigma)}^2.
\]

(2.89)

This shows that \( F_0,\Sigma,(\cdot, 0) \in X^s(R_{\|k\|_{\Sigma}}; F) \), and similarly we can conclude that \( G_0,\Sigma,(\cdot, 0) \in X^s(R_{\|k\|_{\Sigma}}; F) \). Since \( s > d/2 \geq |\Sigma|/2, X^s(R_{\|k\|_{\Sigma}}; F) \) is an algebra, and therefore we have the estimate

\[
\|F_0,\Sigma,(\cdot, 0)G_0,\Sigma,(\cdot, 0)\|_{X^r(R_{\|k\|_{\Sigma}})} \lesssim \|F_0,\Sigma,(\cdot, 0)\|_{X^r(R_{\|k\|_{\Sigma}})}\|G_0,\Sigma,(\cdot, 0)\|_{X^r(R_{\|k\|_{\Sigma}})}.
\]

(2.91)

Combining (2.90) and (2.91) then gives us (2.86) and in turn (2.84). The desired conclusion then follows. \( \square \)

3. LINEAR ANALYSIS

In this section we record the linear analysis associated to the flattened system (1.24). First, we introduce the definition of the \( \hat{H}^{-1} \) seminorm.

**Definition 3.1** (The \( \hat{H}^{-1} \) seminorm). For \( s \geq 0 \), in the case when \( R_\Gamma = \emptyset \), for \( f \in H^s(\Gamma; \mathbb{R}) \) we define the \( \hat{H}^{-1} \) seminorm of \( f \) via

\[
[f]_{\hat{H}^{-1}} = \begin{cases}
0, & \text{if } \hat{f}(0) = 0 \\
\infty, & \text{otherwise}.
\end{cases}
\]

(3.1)

In the remaining cases, we define the \( \hat{H}^{-1} \) seminorm of \( f \) via

\[
[f]_{\hat{H}^{-1}} = \left( \int_{B_{\|k\|_{\Sigma}}(0, r)} |\omega|^{-2} |\hat{f}(\mathcal{P}_\Gamma^{-1} l_{\Gamma} \omega)|^2 d\omega \right)^{1/2},
\]

(3.2)

where \( \Gamma_* \) is defined in (A.8), \( \mathcal{P}_\Gamma \cdot \Gamma \to \Gamma_* \) is defined in (A.10), and \( l_{\Gamma} : \mathbb{R}^{\|\Gamma\|} \to \Gamma_* \) is defined via \( l_{\Gamma}(\omega) = (\omega, 0) \).

For \( f \in H^s(\Gamma; \mathbb{R}) \), we write \( f \in \hat{H}^{-1}(\Gamma; \mathbb{R}) \) to mean that \( [f]_{\hat{H}^{-1}} < \infty \).
We note that our definition of the $H^{-1}$ seminorm coincides with the standard definition of the $H^{-1}$ seminorm for functions $H^s(\Gamma; \mathbb{R})$. First, if a function $f$ already belongs to $H^s(\Gamma; \mathbb{R})$, then

$$
\int_{\hat{\Gamma}} |\xi|^{-2} |\hat{f}(\xi)|^2 \, d\xi \approx \int_{B_{\hat{\Gamma}}(0,r)} |\xi|^{-2} |\hat{f}(\xi)|^2 \, d\xi. \quad (3.3)
$$

Second, we note that the product measure on $\hat{\Gamma}$ is a product of Lebesgue measures and counting measures, corresponding to the $\mathbb{R}$ and $L_i^{-1}\mathbb{Z}$ factors of $\hat{\Gamma}$. If $R_{\hat{\Gamma}} = \emptyset$, then the only low frequency mode is the $0$ mode, and thus the integral in (3.3) takes values in $\{0, \infty\}$ depending on the value of $\hat{f}(0)$. If $R_{\hat{\Gamma}} \neq \emptyset$, then we note that by Lemma A.6, using the natural isometric and measure-preserving identification between $B_{\hat{\Gamma}}(0,r)$ and $B_{\mathbb{R}/n_R}(0,r)$ we have

$$
\int_{B_{\hat{\Gamma}}(0,r)} |\xi|^{-2} |\hat{f}(\xi)|^2 \, d\xi \approx \int_{B_{\mathbb{R}/n_R}(0,r)} |\omega|^{-2} |\hat{f}(\mathcal{P}_{\hat{\Gamma}}^{-1} l_{\hat{\Gamma}} \omega)|^2 \, d\omega. \quad (3.4)
$$

### 3.1. Asymptotics of the normal stress to solution map associated to (1.37)

In this subsection we record the asymptotics of some special functions associated to the normal stress to solution map associated to (1.37), which will play a crucial role in the construction of the free surface function $\eta$. In the case where $\hat{\Sigma} = \mathbb{R}^{n-1}$, the precise asymptotics of these symbols were derived using ODE techniques developed in [18]. For the other possibilities of $\Sigma$ considered in (1.1), these estimates will continue to hold as $\hat{\Sigma}$ remains a subgroup of $\mathbb{R}^{n-1}$ and the asymptotics remain the same upon subsampling.

First we record an auxiliary result.

**Theorem 3.3** (Hilbert isomorphism for the $\gamma$-Stokes system). Suppose $s \geq 0$. The map

$$
\Phi_\gamma : H^{s+1/2}(\Omega; \mathbb{R}^n) \times H^{s+1}(\Omega; \mathbb{R}) \to H^s(\Omega; \mathbb{R}^n) \times H^{s+1}(\Omega; \mathbb{R}) \times H^{s+1/2}(\Sigma_b; \mathbb{R}^n) \quad (3.5)
$$

defined via

$$
\Phi_\gamma(u, p) = (\text{div } S(p, u) - \gamma \partial_1 u, \text{div } u, S(p, u)e_n) \quad (3.6)
$$

is a Hilbert isomorphism for all $\gamma \in \mathbb{R}$.

**Proof.** The case $\Sigma = \mathbb{R}^{n-1}$ follows from Theorem 2.6 in [18], and the arguments therein can be adapted with minimal modification to handle the cases when $\Sigma \neq \mathbb{R}^{n-1}$.  

Next, we define the normal stress to velocity and pressure maps induced by $\Phi_\gamma$.

**Definition 3.4** (Normal stress to velocity and pressure maps). Let $\gamma \in \mathbb{R}$ and $s \geq 0$. We define the normal stress to velocity and pressure maps to be the linear maps $\mathcal{U}_\gamma : H^{s+1/2}(\Omega; \mathbb{R}) \to H^{s+2}(\Omega; \mathbb{R}^n)$ and $\mathcal{P}_\gamma : H^{s+1/2}(\Omega; \mathbb{R}) \to H^{s+1}(\Omega; \mathbb{R})$ defined via

$$
\mathcal{U}_\gamma(\psi) = u, \quad \mathcal{P}_\gamma(\psi) = p, \quad (3.7)
$$

where $(u, p)$ is the unique solution to the $-\gamma$-Stokes system

$$
\begin{align*}
\text{div } S(p, u) + \gamma \partial_1 u &= 0, & \text{in } \Omega \\
\text{div } u &= 0, & \text{in } \Omega \\
S(p, u)e_n &= \psi e_n, & \text{on } \Sigma_b \\
u &= 0, & \text{on } \Sigma_0.
\end{align*} \quad (3.8)
$$

The existence and uniqueness of $(u, p)$ and the boundedness of $\mathcal{U}_\gamma, \mathcal{P}_\gamma$ are guaranteed by Theorem 3.3.

Now we may record properties of the pseudodifferential operator associated to the normal stress to solution maps (3.7).

**Theorem 3.5** (Symbols associated to the pseudodifferential operator and their asymptotics). Suppose $s \geq 0$. The linear maps $\mathcal{U}_\gamma, \mathcal{P}_\gamma$ defined in (3.7) are well-defined and bounded. Moreover, there exists bounded and measurable functions $V : \Sigma \times [0, b] \times \mathbb{R} \to \mathbb{C}^n, Q : \hat{\Sigma} \times [0, b] \times \mathbb{R} \to \mathbb{C}^n$, and $m : \Sigma \times \mathbb{R} \to \mathbb{C}$ such that for all $\psi \in H^{s+1/2}(\Sigma_b; \mathbb{R})$, we have

$$
\mathcal{U}_\gamma(\psi)(\xi, x_n) = V(\xi, x_n, -\gamma) \hat{\psi}(\xi), \quad \mathcal{P}_\gamma(\psi)(\xi, x_n) = Q(\xi, x_n, -\gamma) \hat{\psi}(\xi), \quad m(\xi, \gamma) = V_n(\xi, b, \gamma). \quad (3.9)
$$

Furthermore, the following hold.

1. $V, Q, m$ are continuous, with $V(0, x_n, \gamma) = 0, Q(0, x_n, \gamma) = 1, \text{ and } m(0, \gamma) = 0$.
2. $\mathcal{V}(\xi, x_n, \gamma) = V(-\xi, x_n, \gamma), \mathcal{Q}(\xi, x_n, \gamma) = Q(-\xi, x_n, \gamma), m(\xi, \gamma) = m(-\xi, \gamma)$ for all $\xi \in \hat{\Sigma}$.
(3) For each $\xi \in \hat{\Sigma}$, $V(\xi, \cdot, \gamma), Q(\xi, \cdot, \gamma)$ solve

$$
\begin{align*}
&\left( -\partial_n^2 + 4\pi^2 |\xi|^2 \right) V' + 2\pi i \xi Q - 2\pi i \xi_1 \gamma V' = 0, \quad \text{in } (0, b) \\
&\left( -\partial_n^2 + 4\pi^2 |\xi|^2 \right) V_n + \partial_n Q - 2\pi i \xi_1 \gamma V_n = 0, \quad \text{in } (0, b) \\
&2\pi i \xi \cdot V' + \partial_n V_n = 0, \quad \text{in } (0, b) \\
&-\partial_n V' - 2\pi i \xi V_n = 0, \quad \text{for } x_n = b \\
&Q - 2\partial_n V_n = 1, \quad \text{for } x_n = b \\
&V = 0, \quad \text{for } x_n = 0.
\end{align*}
$$

(3.10)

(4) For $|\xi| \ll 1$, we have

$$
V'(\xi, x_n, \gamma) = -\pi \xi (2x_n b - x_n^2) + O(|\xi|^3), \quad V_n(\xi, x_n, \gamma) = 2\pi^2 |\xi|^2 x_n^2 \left( \frac{x_n^2}{3} - b \right) + O(|\xi|^3),
$$

$$
Q(\xi, x_n, \gamma) = 1 + O(|\xi|^2), \quad m(\xi, \gamma) = \frac{-4\pi^2 |\xi|^2 b^3}{3} + O(|\xi|^3), \quad \text{as } |\xi| \to 0
$$

(3.11)

where $F(\xi, x_n) = O(|\xi|^k)$ as $|\xi| \to 0$ means that

$$
\limsup_{|\xi| \to 0} \sup_{0 \leq x_n \leq b} \frac{|F(\xi, x_n)|}{|\xi|^k} < \infty.
$$

(3.12)

(5) For each $\gamma \in \mathbb{R}$, there exists a constant $c = c(\gamma, b) > 0$ and $R = R(\gamma, b) > 0$ such that for $x_n \in [0, b]$ and $|\xi| > R$, we have the point-wise estimates

$$
|V'(\xi, x_n, \gamma)| \leq c \left( \frac{|\gamma|}{|\xi|^2} + (b - x_n) \right) e^{-2\pi |\xi|(b - x_n)} + ce^{-2\pi |\xi|^b},
$$

(3.13)

$$
|V_n(\xi, x_n, \gamma)| \leq c \left( \frac{1}{|\xi|} + (b - x_n) \right) e^{-2\pi |\xi|(b - x_n)} + ce^{-2\pi |\xi|^b},
$$

(3.14)

$$
|m(\xi, \gamma) + \frac{1}{4\pi |\xi|}| \leq c \frac{1}{|\xi|^2},
$$

(3.15)

$$
|Q(\xi, x_n, \gamma)| \leq ce^{-2\pi |\xi|(b - x_n)} + ce^{-2\pi |\xi|^b}.
$$

(3.16)

(6) For each $\gamma \in \mathbb{R}$, there exists a constant $c = c(n, b, \gamma)$ such that

$$
\sup_{|\xi| \leq 1, \xi \neq 0} \frac{1}{|\xi|^2} \left( \int_0^b |V(\xi, x_n, \gamma)|^2 + |Q(\xi, x_n, \gamma) - 1|^2 \, dx_n + |V(\xi, b, \gamma)|^2 \right) \leq c
$$

(3.17)

and

$$
(1 + |\xi|^2)^{3/2} \int_0^b |V(\xi, x_n, \gamma)|^2 \, dx_n + (1 + |\xi|^2)^{1/2} \int_0^b |Q(\xi, x_n, \gamma)|^2 \, dx_n + (1 + |\xi|^2)|V(\xi, b, \gamma)|^2 \leq c,
$$

(3.18)

for all $\xi \in \hat{\Sigma}$.

Proof. If $\hat{\Sigma} = \mathbb{R}^{n-1}$, all six items follow from Theorem 4.5, Theorem 4.10, and Corollary 4.11 in [18]. In the other cases, we have $\hat{\Sigma} \subset \mathbb{R}^{n-1} = \mathbb{R}^{n-1}$. Therefore, the same conclusions will follow as we are sampling the frequencies on subgroups of $\mathbb{R}^{n-1}$.

We conclude this subsection by recording the properties of an auxiliary function defined in terms of $m$.

Lemma 3.6. Suppose $\gamma \in \mathbb{R} \setminus \{0\}$, and define

$$
\rho_\gamma(\xi) = 2\pi i \gamma \xi_1 + (1 + 4\pi^2 |\xi|^2) s m(\xi, -\gamma).
$$

(3.19)

Then the following hold.

(1) $\rho_\gamma(\xi) = 0$ if and only if $\xi = 0$, and $\rho_\gamma(-\xi) = \rho_\gamma(\xi)$ for all $\xi \in \hat{\Sigma}$.

(2) For $\sigma > 0$, there exists a constant $C = C(n, \gamma, \sigma, b) > 0$ such that for all $\xi \in \hat{\Sigma}$, we have

$$
C^{-1} |\rho_\gamma(\xi)|^2 \leq (\xi_1^2 + |\xi|^4) \mathbb{I}_{B(0,1)}(\xi) + (1 + |\xi|^2) \mathbb{I}_{B(0,1)^c}(\xi) \leq C |\rho_\gamma(\xi)|^2.
$$

(3.20)

(3) For $\sigma = 0$ and $n = 2$, there exists a constant $C = C(\gamma, b) > 0$ such that for all $\xi \in \hat{\Sigma}$, we have

$$
C^{-1} |\rho_\gamma(\xi)|^2 \leq (|\xi|^2 \mathbb{I}_{B(0,1)}(\xi) + (1 + |\xi|^2) \mathbb{I}_{B(0,1)^c}(\xi) \leq C |\rho_\gamma(\xi)|^2.
$$

(3.21)

Proof. This follows from the proof for Lemma 6.1 in [18] and the asymptotics of $m(\xi, -\gamma)$ on general $\hat{\Sigma}$ recorded as the third and fourth items of Theorem 3.5.
3.2. Compatibility conditions and the Hilbert isomorphism associated to the overdetermined \( \gamma \)-Stokes system. In this subsection we record the Hilbert isomorphism that completely characterize the solvability of the overdetermined system (1.37). First we introduce a pair of function spaces for the data tuple \((f, g, h, k)\) that encodes the compatibility conditions associated to (1.37).

**Definition 3.7.** Let \( s \geq 0 \) and \( \Sigma \) be defined as in (1.1). We say that a data tuple \((f, g, h, k) \in H^s(\Omega; \mathbb{R}^n) \times H^{s+1}(\Omega; \mathbb{R}) \times H^{s+1/2}(\Sigma; \mathbb{R}) \times H^{s+1/2}(\Sigma; \mathbb{R}^n)\) satisfy the divergence-trace compatibility condition if

\[
h - \int_0^b g(\cdot, x_n) \, dx_n \in H^{-1}(\Sigma; \mathbb{R}).
\]

We define the Hilbert space \( \mathcal{Y}^s \) to be

\[
\mathcal{Y}^s = \{(f, g, h, k) \in H^s(\Omega; \mathbb{R}^n) \times H^{s+1}(\Omega; \mathbb{R}) \times H^{s+3/2}(\Sigma; \mathbb{R}) \times H^{s+1/2}(\Sigma; \mathbb{R}^n) \mid (f, g, h, k) \text{ satisfy the divergence-trace condition (3.22)}\},
\]

with the associated norm defined via

\[
\|(f, g, h, k)\|_{\mathcal{Y}^s}^2 = \|f\|^2_{H^s} + \|g\|^2_{H^{s+1}} + \|h\|^2_{H^{s+3/2}} + \left[ h - \int_0^b g(\cdot, x_n) \, dx_n \right]_{H^{-1}}^2.
\]

**Definition 3.8.** Let \( s \geq 0 \). We say that the data tuple \((f, g, h, k) \in H^s(\Omega; \mathbb{R}^n) \times H^{s+1}(\Omega; \mathbb{R}) \times H^{s+3/2}(\Sigma; \mathbb{R}) \times H^{s+1/2}(\Sigma; \mathbb{R}^n)\) satisfy the adjoint compatibility condition if for every \( \psi \in H^{s+1/2}(\Sigma; \mathbb{R})\),

\[
\int_0^b \hat{f}(\xi, x_n) \cdot \hat{V}(\xi, x_n, -\gamma) - \hat{g}(\xi, x_n) \hat{Q}(\xi, x_n, -\gamma) \, dx_n - \hat{h}(\xi) \cdot \hat{V}(\xi, 0, -\gamma) + \hat{h}(\xi) = 0,
\]

where \( V, Q \) are defined in terms of \( \psi \) in (3.10). We define the closed subspace \( \mathcal{Z}^s \) of \( \mathcal{Y}^s \) to be

\[
\mathcal{Z}^s = \{(f, g, h, k) \in \mathcal{Y}^s \mid (f, g, h, k) \text{ satisfy (3.25)} \forall \psi \in H^{s+1/2}(\Sigma; \mathbb{R})\}.
\]

Now we are ready to record the main result of this subsection.

**Theorem 3.9.** Let \( s \geq 0, \gamma \in \mathbb{R}, \) and \( \mathcal{Z}^s \) the Hilbert space defined in (3.8). Then the bounded linear operator

\[
\Psi_\gamma : 0H^{s+2}(\Omega; \mathbb{R}^n) \times H^{s+1}(\Omega; \mathbb{R}) \to \mathcal{Z}^s \text{ given by }
\]

\[
\Psi_\gamma(u, p) = (\text{div } S(p, u) - \gamma \partial_1 u, \text{div } u, u_n|_{\Sigma_b}, S(p, u) e_n|_{\Sigma_b})
\]

is an isomorphism.

**Proof.** The case \( \Sigma = \mathbb{R}^{n-1} \) follows from Theorem 3.4 of [18]. We note that by Remark 3.2, the calculations performed in \( \mathbb{R}^d \) in [18] are also valid over the low frequency regime in \( \Gamma \), therefore the arguments therein can be adapted with minimal modification to handle the cases when \( \Sigma \neq \mathbb{R}^{n-1} \).

3.3. Linear analysis with \( \eta \) and \( \kappa \). In this section we would like to establish the solvability of the \( \gamma \)-Stokes system with gravity capillary boundary conditions

\[
\begin{align*}
\text{div } S(p, u) - \gamma \partial_1 u + (\nabla \eta, 0) &= f \quad \text{in } \Omega, \\
\text{div } u &= g, \quad \text{in } \Omega, \\
u_n + \gamma \partial_1 \eta &= h, \quad \text{on } \Sigma_b, \\
S(p, u) e_n + \sigma \Delta' \eta e_n &= k, \quad \text{on } \Sigma_b, \\
u &= 0, \quad \text{on } \Sigma_0,
\end{align*}
\]

for data tuples belonging to the space \( \mathcal{Y}^s \). First, we introduce the container space for the solution tuple \((u, p, \eta)\).

**Definition 3.10.** For \( s \geq 0 \), we define the Hilbert space

\[
\mathcal{X}^s = \begin{cases} 
\{(u, p, \eta) \in 0H^{s+2}(\Omega; \mathbb{R}^n) \times H^{s+1}(\Omega; \mathbb{R}) \times X^{s+5/2}(\Sigma; \mathbb{R}) \mid \eta \in X^{s+5/2}(\Sigma; \mathbb{R}) \} & R_\Sigma \neq \emptyset \\
\{(u, p, \eta) \in 0H^{s+2}(\Omega; \mathbb{R}^n) \times H^{s+1}(\Omega; \mathbb{R}) \times \hat{X}^{s+5/2}(\Sigma; \mathbb{R}) \} & R_\Sigma = \emptyset,
\end{cases}
\]

where

\[
\hat{X}^{s+5/2}(\Sigma; \mathbb{R}) = \{ \eta \in X^{s+5/2}(\Sigma; \mathbb{R}) \mid \eta(0) = 0 \}.
\]

We endow \( \mathcal{X}^s \) with the natural product norm defined via

\[
\|(u, p, \eta)\|_{\mathcal{X}^s}^2 = \|u\|_{0H^{s+1}}^2 + \|p\|_{H^{s+1}}^2 + \|\eta\|_{\hat{X}^{s+5/2}}^2.
\]

Next, we record an embedding result for \( \mathcal{X}^s \).
Proposition 3.11. Suppose $s \geq 0$ and $\mathcal{X}^s$ is the Banach space in Definition 3.10. If $s > n/2$, then we have the continuous inclusion
\[
\mathcal{X}^s \subseteq C_b^{[s-n/2]}(\Omega; \mathbb{R}^n) \times C_b^{1+[s-n/2]}(\Omega; \mathbb{R}) \times C_b^{3+[s-n/2]}(\Sigma; \mathbb{R}).
\] (3.32)

Moreover, if $(u,p,\eta) \in \mathcal{X}^s$, then
\[
\lim_{|x'_{n}| \to \infty} \partial^\alpha u(x) = 0 \text{ for all } \alpha \in \mathbb{N}^n \text{ such that } |\alpha| \leq 2 + \left[s - \frac{n}{2}\right]
\] (3.33)
\[
\lim_{|x'_{n}| \to \infty} \partial^\alpha p(x) = 0 \text{ for all } \alpha \in \mathbb{N}^n \text{ such that } |\alpha| \leq 1 + \left[s - \frac{n}{2}\right].
\] (3.34)

Proof. This follows from standard Sobolev embedding and the first item of Theorem 2.6.

In the subsection to follow, we establish some preliminary results to be utilized in the subsequent analysis.

3.4. Preliminary results. In this subsection we use the asymptotics recorded in Section 3.1 to show that we can construct the free surface function $\eta$ from a given data tuple $(f,g,h,k)$. First, we study an auxiliary function defined in terms of the multipliers $V,Q$ defined in (3.9).

Lemma 3.12. Let $s \geq 0$, $\gamma \in \mathbb{R}$ and $(f,g,h,k) \in \mathcal{Y}^s$, where $\mathcal{Y}^s$ is the Hilbert space defined in (3.23). Consider $V(\cdot,\cdot,\gamma) : \hat{\Sigma} \times [0,b] \to \mathbb{C}^n$, $Q(\cdot,\cdot,\gamma) : \hat{\Sigma} \times [0,b] \to \mathbb{C}$ defined in (3.9). We define the measurable function $\psi : \hat{\Sigma} \to \mathbb{C}$ via
\[
\psi(\xi) = \int_0^b \left( \hat{f}(\xi,x_n) \cdot \overline{V(\xi,x_n,-\gamma)} - \hat{g}(\xi,x_n)Q(\xi,x_n,-\gamma) \right) \; dx_n - \hat{k}(\xi) \cdot \overline{V(\xi,b,-\gamma)} + \hat{h}(\xi).
\] (3.35)

Then the following hold.
1. If $R_\Sigma = \emptyset$, then $\psi(0) = 0$.
2. $\psi(\xi) = \psi(-\xi)$ for every $\xi \in \hat{\Sigma}$.
3. We have the estimate
\[
\int_{\hat{\Sigma}} \frac{1}{|\xi|^2} |\psi(\xi)|^2 1_{B(0,r)} + (1 + |\xi|^2)^{s+3/2} |\psi(\xi)|^2 1_{B(0,\gamma)} \; d\xi \lesssim_{n,s,\eta,b} \|(f,g,h,k)\|^2_{\mathcal{Y}^s}.
\] (3.36)

Proof. To prove the first item, we assume $R_\Sigma = \emptyset$ and rewrite
\[
\psi(\xi) = \int_0^b \left( \hat{f}(\xi,x_n) \cdot \overline{V(\xi,x_n,-\gamma)} - \hat{g}(\xi,x_n)(Q(\xi,x_n,-\gamma) - 1) \right) \; dx_n
\]
\[
- \hat{k}(\xi) \cdot \overline{V(\xi,b,-\gamma)} + \left( \hat{h}(\xi) - \int_0^b \hat{g}(\xi,x_n) \; dx_n \right).
\] (3.37)

We note that since $h,g$ must satisfy (3.22), by Remark 3.2 we have $\hat{h}(0) - \int_0^b \hat{g}(0,x_n) \; dx_n = 0$. We also note that by the first item of Theorem 3.5, $\overline{V(0,x_n,-\gamma)} = 0$ and $Q(\xi,x_n,-\gamma) = 0$. The first item follows immediately from these two observations.

To prove the second item, we first note that by using the second item of Theorem 3.5 and by using the fact that $f,g,h,k$ are real-valued,
\[
\overline{\psi}(\xi) = \int_0^b \left( \hat{f}(\xi,x_n) \cdot V(-\xi,x_n,-\gamma) - \hat{g}(\xi,x_n)Q(-\xi,x_n,-\gamma) \right) \; dx_n + \hat{k}(-\xi) \cdot V(-\xi,b,-\gamma) + \hat{h}(-\xi)
\]
\[
= \int_0^b \left( \hat{f}(-\xi,x_n) \cdot V(-\xi,x_n,-\gamma) - \hat{g}(-\xi,x_n)Q(-\xi,x_n,-\gamma) \right) \; dx_n - \hat{k}(-\xi) \cdot \overline{V(-\xi,b,-\gamma)} + \hat{h}(-\xi) = \psi(-\xi).
\] (3.38)

To prove the third, we first rewrite $\psi$ again as in (3.37) and apply the Cauchy-Schwarz inequality to obtain
\[
|\psi(\xi)|^2 \lesssim \left( \int_0^b |\hat{f}(\xi,x_n)|^2 \; dx_n \right) \left( \int_0^b |V(\xi,x_n)|^2 \; dx_n \right) + \left( \int_0^b |\hat{g}(\xi,x_n)|^2 \; dx_n \right) \left( \int_0^b |Q(\xi,x_n,-\gamma) - 1|^2 \; dx_n \right)
\]
\[
+ |\hat{k}(\xi)|^2 |V(\xi,b,-\gamma)|^2 + |\hat{h}(\xi)| - \int_0^b |\hat{g}(\xi,x_n) |^2 \; dx_n^2.
\] (3.39)
We note that by the definition of the $\dot{H}^{-1}$ seminorm, we have

$$\int_{B(0,r)} \frac{1}{|x|^2} \hat{h}(\xi) - \int_0^b \hat{g}(\xi, x_n) \, dx_n^2 \, d\xi \leq \left[ h - \int_0^b g(\cdot, x_n) \, dx_n \right]_{\dot{H}^{-1}}^2. \quad (3.40)$$

Combining (3.39), (3.40) with (3.17), Tonelli’s theorem, and Parseval’s theorem then gives us

$$\int_{B(0,r)} \frac{1}{|\xi|^2} |\psi(\xi)|^2 \, d\xi \lesssim_{n,\gamma,b} \int_{B(0,r)} \int_0^b \left( |\hat{f}(\xi, x_n)|^2 + |\hat{g}(\xi, x_n)|^2 \right) \, dx_n \, d\xi$$

$$+ \int_{B(0,r)} |\hat{k}(\xi)|^2 \, d\xi + \left[ h - \int_0^b g(\cdot, x_n) \, dx_n \right]_{\dot{H}^{-1}}^2 \lesssim_{n,\gamma,b} \| (f, g, h, k) \|^2_{Y^s}. \quad (3.41)$$

If $|\xi| \geq 1$, then we apply the Cauchy-Schwarz inequality directly to (3.35) and obtain

$$|\psi(\xi)|^2 \lesssim \left( \int_0^b |\hat{f}(\xi, x_n)|^2 \, d\xi \right) \left( \int_0^b |V(\xi, x_n)|^2 \, d\xi \right) + \left( \int_0^b |\hat{g}(\xi, x_n)|^2 \, d\xi \right) \left( \int_0^b |Q(\xi, x_n, -\gamma)|^2 \, d\xi \right)$$

$$+ |\hat{k}(\xi)|^2 |V(\xi, b, -\gamma)|^2 + |\hat{h}(\xi)|^2. \quad (3.42)$$

Combining (3.42) with (3.18), Tonelli’s theorem, and Parseval’s theorem then gives us

$$\int_{B(0,r)} (1 + |\xi|^2)^{s+3/2} |\psi(\xi)|^2 \, d\xi \lesssim_{n,\gamma,b} \int_{B(0,r)} (1 + |\xi|^2)^s |\hat{f}(\xi, x_n)|^2 \, d\xi \, dx_n + \int_0^b \int_{\Sigma} (1 + |\xi|^2)^s |\hat{g}(\xi, x_n)|^2 \, d\xi \, dx_n$$

$$+ \int_{\Sigma} (1 + |\xi|^2)^{s+1/2} |\hat{k}(\xi)|^2 \, d\xi + \int_{\Sigma} (1 + |\xi|^2)^{s+3/2} |\hat{h}(\xi)|^2 \, d\xi \lesssim_{n,\gamma,b} \int_0^b \| f(\cdot, x_n) \|^2_{H^{s}(\Sigma)} \, dx_n$$

$$+ \int_0^b \| g(\cdot, x_n) \|^2_{H^{s+1}(\Sigma)} \, dx_n + \| k \|^2_{H^{s+1/2}} + \| h \|^2_{H^{s+3/2}} \lesssim_{n,\gamma,b} \| (f, g, h, k) \|^2_{Y^s}. \quad (3.43)$$

Next we study the linear map $\Upsilon_{\gamma,\sigma} : X^s \to Y^s$ defined via

$$\Upsilon_{\gamma,\sigma}(u, p, \eta) = (\text{div} S(p, u) - \gamma \partial_1 u + (\nabla \eta, 0), \text{div} u, u, u_n|_{\Sigma_b} + \gamma \partial_1 \eta, S(p, u) e_n|_{\Sigma_b} + \sigma \Delta \eta e_n) \quad (3.44)$$

which is the solution operator corresponding to the system (3.28). The next result shows that this linear map is well-defined, bounded, and also injective.

**Proposition 3.13.** Suppose $\gamma \in \mathbb{R} \setminus \{0\}$, $\sigma \geq 0$, and $s \geq 0$. Then the linear map $\Upsilon_{\gamma,\sigma} : X^s \to Y^s$ defined in (3.44) is well-defined, continuous, and injective.

**Proof.** We first check that the map is well-defined and continuous. By the fifth item of Theorem 2.5, the first component $f = \text{div} S(p, u) - \gamma \partial_1 u + (\nabla \eta, 0)$ belongs to $H^{s+1}(\Omega; \mathbb{R}^n)$ with the estimate

$$\| f \|_{H^{s+1}} \lesssim \| u \|_{H^{s+2}} + \| p \|_{H^{s+1}} + \| \eta \|_{X^{s+5/2}}. \quad (3.45)$$

The second component $g = \text{div} u$ belongs to $H^{s+1}(\Omega; \mathbb{R})$ with the estimate $\| g \|_{H^{s+1}} \lesssim \| u \|_{H^{s+2}}$. By the fifth item of Theorem 2.5 and standard trace theory, we deduce that the third component $k = S(p, u) e_n|_{\Sigma_b} + \sigma \Delta \eta e_n$ belongs to $H^{s+1}(\Sigma_b; \mathbb{R}^n)$ with the estimate

$$\| k \|_{H^{s+1/2}(\Sigma)} \lesssim \| p \|_{H^{s+1/2}(\Sigma)} + \| D u \|_{H^{s+1/2}(\Sigma)} + \| \sigma \Delta \eta \|_{H^{s+1/2}(\Sigma)} \lesssim \| p \|_{H^{s+1}(\Omega)} + \| u \|_{H^{s+2}} + \| \eta \|_{X^{s+5/2}}. \quad (3.46)$$

By the fifth item of Theorem 2.5 and standard trace theory, the fourth component $h = u_n|_{\Sigma_b} + \gamma \partial_1 \eta$ belongs to $H^{s+3/2}(\Sigma_b; \mathbb{R})$ with the estimate

$$\| h \|_{H^{s+3/2}} \lesssim \| u \|_{H^{s+2}} + \| \eta \|_{X^{s+5/2}}. \quad (3.47)$$

Next we note that since $g = \text{div} u$,

$$u_n(x', b) - \int_0^b g(x', x_n) \, dx_n = \int_0^b \partial_n u_n(x', x_n) \, dx_n - \int_0^b g(x', x_n) = - \text{div} \int_0^b u'(x', x_n) \, dx_n \quad (3.48)$$
for a.e. \( x' \in \Sigma \). Writing \( H^{s+3}(\Sigma; \mathbb{R}) \ni R(x') = \int_0^b u'(x', x_n) \, dx_n \), we first note that if \( R_\Gamma = \emptyset \), then \( \text{div} \, R(0) = 0 \) and \( \gamma \partial_1 \eta(0) = 0 \). Then we note that

\[
\left[ u_n - \int_0^b g(\cdot, x_n) \, dx_n \right]_{\hat{H}^{-1}} = \left[ \text{div} \, R \right]_{\hat{H}^{-1}} \lesssim \| u \|_{L^2},
\]

therefore

\[
\left[ h - \int_0^b g(\cdot, x_n) \, dx_n \right]_{\hat{H}^{-1}} \lesssim \left[ u_n - \int_0^b g(\cdot, x_n) \, dx_n \right]_{\hat{H}^{-1}} + [\gamma \partial_1 \eta]_{\hat{H}^{-1}} \lesssim \| u \|_{aH^{s+2}} + \| \eta \|_{X^{s+5/2}}.
\]

Combining the estimates above shows that \( \Upsilon_{\gamma, \sigma} \) is well-defined and continuous.

To show that \( \Upsilon_{\gamma, \sigma} \) is injective, we suppose \((u, p, \eta) \in X^s \) and \( \Upsilon_{\gamma, \sigma}(u, p, \eta) = 0 \). We note that if \( \tilde{p} = p - \eta \), then \( \nabla \tilde{p} = \nabla p - (\nabla \eta, 0) \) and \( \tilde{p} I = p I - \eta I \). Therefore \( \Upsilon_{\gamma, \sigma}(u, p, \eta) = 0 \) if and only if \((u, \tilde{p}, \eta)\) satisfies

\[
\begin{align*}
\text{div} \, S(\tilde{p}, u) - \gamma \partial_1 u &= 0, & \text{in } \Omega \\
u &= 0, & \text{in } \Omega \\
S(\tilde{p}, u) e_n &= (\eta - \sigma \Delta \eta) e_n, & \text{on } \Omega_b \\
u_n + \gamma \partial_1 \eta &= 0, & \text{on } \Omega_b \\
u &= 0, & \text{on } \Sigma_0.
\end{align*}
\]

We note that by Tonelli’s theorem, Parseval’s theorem, and the fifth item of Theorem 2.5 we have \( \hat{u}(\xi, \cdot) \in H^s((0, b); \mathbb{C}^n) \) and \( \hat{\rho}(\xi, \cdot) \in H^1((0, b); \mathbb{C}) \), for a.e. \( \xi \in \hat{\Sigma} \). By the second item in Theorem 2.6, \( \hat{\eta} \in L^1(\Sigma; \mathbb{R}) + L^2(\Sigma, (1 + |\xi|^2)^{(s+5/2)/2} |d\xi|; \mathbb{R}) \). Thus, we may apply the horizontal Fourier transform to (3.51) to deduce for a.e. \( \xi \in \hat{\Sigma}, w = \hat{u}(\xi, \cdot), q = \hat{\rho}(\xi, \cdot) \) satisfies

\[
\begin{align*}
(-\partial_n^2 + 4\pi^2|\xi|^2) w' + 2\pi i \xi q - 2\pi i \xi_1 \gamma w' &= 0, & \text{in } (0, b) \\
(-\partial_n^2 + 4\pi^2|\xi|^2) w_n + \partial_n q - 2\pi i \xi_1 \gamma w_n &= 0, & \text{in } (0, b) \\
2\pi i \xi \cdot w' + \partial_n w_n &= 0, & \text{in } (0, b) \\
-\partial_n w' - 2\pi i \xi w_n &= 0, & \text{for } x_n = b \\
q - 2\partial_n w_n &= (1 + 4\pi^2|\xi|^2 \sigma) \hat{\eta}, & \text{for } x_n = b \\
w_n + 2\pi i \xi_1 \gamma \hat{\eta} &= 0, & \text{for } x_n = b \\
w &= 0, & \text{for } x_n = 0.
\end{align*}
\] (3.52)

First consider the special case when \( R_\Sigma = \emptyset \) and \( \xi = 0 \). In this case, the third and sixth equations imply that \( w_n \equiv 0 \). Since \( \eta \in \hat{X}^{s+5/2}(\Sigma; \mathbb{R}) \), we have \( \hat{\eta}(0) = 0 \) and therefore by the second and fifth equations, \( q \equiv 0 \). The first, fourth, and last equations then tell us that \( w' \equiv 0 \). Therefore, we can conclude that at \( \xi = 0, w, q \equiv 0 \).

Next consider the general case for which \( \xi \in \hat{\Sigma} \setminus \{0\} \) and (3.52) holds. Using (3.52) and performing integration by parts (for further details, see Proposition 4.1 in [18]), we deduce that

\[
\int_0^b -\gamma 2\pi i \xi_1 |w|^2 + 2|\partial_n w_n|^2 + |\partial_n w' + 2\pi i \xi w_n|^2 + \frac{1}{2} |2\pi i \xi \otimes w' + w' \otimes 2\pi i \xi|^2 \, dx_n = -2\pi i \xi_1 \gamma (1 + 4\pi^2|\xi|^2 \sigma) |\hat{\eta}(\xi)|^2.
\]

By taking the real part of this expression, we see that we must have for a.e. \( \xi \in \hat{\Sigma}, \partial_n w_n \equiv 0 \) and \( \partial_n w' + 2\pi i \xi w_n \equiv 0 \) in \( (0, b) \), but since \( w(0) = 0 \), we must have \( w \equiv 0 \) in \( [0, b] \). Then by the first equation, we must have \( q \equiv 0 \). By the second to last equation, we find that \( \eta \equiv 0 \). From this we find that \((u, p, \eta) = (0, 0, 0)\), so we can conclude that \( \Upsilon_{\gamma, \sigma} \) is injective.

Next we show that \( \Upsilon_{\gamma, \sigma} \) is surjective. To do so we must construct the free surface function \( \eta \) from a given data tuple \((f, g, h, k) \in \mathcal{Y}^s \). For the reader’s convenience, we record this construction in the next subsection.

3.5. Construction of the free surface function and the isomorphism associated to (3.28).

Lemma 3.14 (Construction of \( \eta \) in the presence of surface tension). Suppose \( \gamma \in \mathbb{R} \setminus \{0\} \), \( \sigma > 0 \) and \( n \geq 2, s \geq 0 \). Then for every \((f, g, h, k) \in \mathcal{Y}^s \), there exists an \( \eta \in X^{s+\frac{5}{2}}(\Sigma; \mathbb{R}) \) for which the modified data tuple

\[
(f - (\nabla \eta, 0), g, h - \gamma \partial_1 \eta, k - \sigma \Delta^2 \eta e_n) \in H^s(\Omega; \mathbb{R}^n) \times H^{s+1}(\Omega; \mathbb{R}) \times H^{s+\frac{5}{2}}(\Sigma_0; \mathbb{R}) \times H^{s+\frac{5}{2}}(\Sigma_0; \mathbb{R}^n)
\]

(3.54)
belongs to the range of $\Upsilon_{\gamma,\sigma}$ defined in (3.44). Moreover, we have the estimate
\begin{equation}
\|\eta\|_{X^{s+\frac{1}{2}}} \lesssim \|(f, g, h, k)\|_{Y^s} .
\end{equation}

Proof. Given $(f, g, h, k) \in Y^s$, we propose to define $\eta \in X^{s+\frac{1}{2}}(\Sigma; \mathbb{R})$ through
\begin{equation}
\hat{\eta}(\xi) = \begin{cases} \frac{\psi(\xi)}{\rho(\xi)}, & \xi \neq 0, \\
0, & \xi = 0,
\end{cases}
\end{equation}
where $\rho_\gamma$ is defined in (3.19) and $\psi$ is defined in terms of $(f, g, h, k)$ in (3.35). We note that the choice of $\hat{\eta}(0) = 0$ is only relevant in the case when $R_\Gamma = \emptyset$.

We first note that by the first item of Lemma 3.6 and the second item of Lemma 3.12, we have $\overline{\eta(\xi)} = \hat{\eta}(-\xi)$. Furthermore, by the second item of Lemma 3.6 and the third item of Lemma 3.12, we have the estimate
\begin{equation}
\int_\Sigma \left( \frac{\xi^2 + |\xi|^4}{|\xi|^2} 1_{B(0, r)} + (1 + |\xi|^2)^{s+\frac{1}{2}} 1_{B(0, r)^c} \right) |\hat{\eta}(\xi)|^2 d\xi 
\lesssim \int_\Sigma \left( \frac{1}{|\xi|^2} 1_{B(0, r)} + (1 + |\xi|^2)^{s+\frac{1}{2}} 1_{B(0, r)^c} \right) |\rho_\gamma(\xi)|^2 |\hat{\eta}(\xi)|^2 d\xi
\lesssim \int_\Sigma \left( \frac{1}{|\xi|^2} 1_{B(0, r)} + (1 + |\xi|^2)^{s+\frac{1}{2}} 1_{B(0, r)^c} \right) |\psi(\xi)|^2 d\xi \lesssim \|(f, g, h, k)\|_{Y^s}^2 .
\end{equation}

Consequently, if we define $\eta = (\hat{\eta})^\vee$, then by (3.57), and $\overline{\eta(\xi)} = \hat{\eta}(\xi)$, we have $\eta \in X^{s+\frac{1}{2}}(\Sigma; \mathbb{R})$ with the estimate (3.55).

To conclude our proof it suffices to show that the modified data given in (3.54) belongs to the range of $\Upsilon_{\gamma,\sigma}(u, p, \eta)$. By Theorem 3.9, it suffices to show that the modified data tuple has the desired regularity and satisfies the divergence-trace compatibility condition (3.22) and the adjoint compatibility condition (3.25). We note that since $\eta \in X^{s+\frac{1}{2}}(\Sigma; \mathbb{R})$, by the fifth item of Theorem 2.5 and the third item of Theorem 2.6, we have $(f, g, h - \gamma \partial_1 \eta, k - \sigma \Delta' \eta \eta_{n}) \in H^s(\Omega; \mathbb{R}^n) \times H^{s+1}(\Omega; \mathbb{R}) \times H^{s+\frac{1}{2}}(\Sigma_b; \mathbb{R}) \times H^{s+\frac{1}{2}}(\Sigma_b; \mathbb{R}^n)$ with $(h - \gamma \partial_1 \eta, - \int_0^b g(\cdot, x_n) \in H^{-1}(\Sigma; \mathbb{R})$. To check (3.25), we write $0 = \psi - \rho_\gamma \hat{\eta}$ for $\xi \neq 0$ and use the definition of $\psi$ and $\rho_\gamma$ to obtain
\begin{equation}
0 = \int_0^b \left( \hat{f}(\xi, x_n) \cdot V(\xi, x_n, -\gamma) - \hat{g}(\xi, x_n) Q(\xi, x_n, -\gamma) \right) dx_n \\
- (\hat{k}(\xi) + (1 + 4\pi^2 |\xi|^2 \sigma)\hat{\eta}(\xi)c_n) \cdot V(\xi, b, -\gamma) + \hat{h}(\xi) - 2\pi i \xi_1 \hat{\eta}(\xi).
\end{equation}

Using the third equation in (3.10) we have $2\pi i \xi \cdot V'(\xi, x_n, -\gamma) + \partial_n V_n(\xi, x_n, -\gamma) = 0$, and since $V(\xi, 0, -\gamma) = 0$, we have
\begin{equation}
V_n(\xi, b, -\gamma) = \int_0^b \partial_n V_n(\xi, x_n, -\gamma) dx_n = \int_0^b 2\pi i \xi \cdot V'(\xi, x_n, -\gamma) dx_n.
\end{equation}

Thus, we can rewrite (3.58) as
\begin{equation}
0 = \int_0^b \left( \left( \hat{f}(\xi, x_n) - ((\nabla' \eta)(\xi), 0) \right) \cdot V(\xi, x_n, -\gamma) - \hat{g}(\xi, x_n) Q(\xi, x_n, -\gamma) \right) dx_n \\
- (\hat{k}(\xi) + \sigma \Delta' \eta(c_n)) \cdot V(\xi, b, -\gamma) + \hat{h}(\xi) - 2\pi i \xi_1 \hat{\eta}(\xi),
\end{equation}
and the desired conclusion follows immediately.

For the special case of $n = 2$, we can also construct the free surface function $\eta$ in the case without surface tension.

Lemma 3.15 (Construction of the free surface function without surface tension). Suppose $\gamma \in \mathbb{R} \setminus \{0\}$, $\sigma = 0$ and $n = 2$, $s \geq 0$. Then for every $(f, g, h, k) \in Y^s$, there exists an $\eta \in H^{s+\frac{1}{2}}(\Sigma; \mathbb{R})$ for which the modified data tuple
\begin{equation}
(f - \partial_1 \eta e_1, g, h - \gamma \partial_1 \eta, k) \in H^s(\Omega; \mathbb{R}^n) \times H^{s+1}(\Omega; \mathbb{R}) \times H^{s+\frac{1}{2}}(\Sigma_b; \mathbb{R}) \times H^{s+\frac{1}{2}}(\Sigma_b; \mathbb{R}^n)
\end{equation}
belongs to the range of $\Upsilon_{\gamma,\sigma}$ defined in (3.44). Moreover, we have the estimate
\begin{equation}
\|\eta\|_{X^{s+\frac{1}{2}}} \lesssim \|(f, g, h, k)\|_{Y^s} .
\end{equation}
Proof. We follow the construction in the previous lemma and propose to define \( \eta \) via (3.56). Then \( \overline{\eta}(\xi) = \hat{\eta}(-\xi) \), \( \hat{\eta}(0) = 0 \), and by the third item of Lemma 3.6 and the third item of Lemma 3.12, we have
\[
\int_{\Sigma} \left( (1 + |\xi|^2)\mathbf{1}_{B(0,1)} + (1 + |\xi|^2)^{s+\frac{3}{2}}\mathbf{1}_{B(0,1)^c} \right) |\hat{\eta}(\xi)|^2 \, d\xi
\leq \int_{\Sigma} \left( \frac{1}{|\xi|^2} \mathbf{1}_{B(0,1)} + (1 + |\xi|^2)^{s+\frac{3}{2}} \mathbf{1}_{B(0,1)^c} \right) |\rho_\gamma(\xi)|^2 |\hat{\eta}(\xi)|^2 \, d\xi
\leq \int_{\Sigma} \left( \frac{1}{|\xi|^2} \mathbf{1}_{B(0,1)} + (1 + |\xi|^2)^{s+\frac{3}{2}} \mathbf{1}_{B(0,1)^c} \right) |\psi(\xi)|^2 \, d\xi
\leq \|(f, g, h, k)\|_{L^p}^2. \tag{3.63}
\]
Consequently, we may define \( \eta = (\hat{\eta})^\vee \). By (3.63) and Lemma 2.2, we have \( \eta \in H^{s+5/2}(\Sigma; \mathbb{R}) = X^{s+5/2}(\Sigma; \mathbb{R}) \) with the estimate (3.62). By the third item of Theorem 2.6, we have \( (f - \partial_1 \eta e_1, g, -\gamma \partial_1 \eta, k) \in \mathcal{H}^s(\Omega; \mathbb{R}^n) \times H^{s+1}(\Omega; \mathbb{R}) \times \mathcal{H}^{s+\frac{3}{2}}(\Sigma_b; \mathbb{R}) \times \mathcal{H}^{s+\frac{3}{2}}(\Sigma_b; \mathbb{R}^n) \) and \( (h - \gamma \partial_1 \eta) - \int_0^\delta g(\cdot, x_2) \, dx_2 \in \mathcal{H}^{-1}(\Sigma; \mathbb{R}) \). To conclude we follow the same line of calculations as in Lemma 3.11 to arrive at
\[
0 = \int_0^\delta \left( (f(\xi, x_2) - 2\pi i \xi_1 \hat{\eta}(\xi_1)) \cdot \nabla(\xi, x_2, -\gamma) - \hat{g}(\xi, x_2) Q(\xi, x_2, -\gamma) \right) \, dx_2 - \hat{h}(\xi) \cdot \nabla(\xi, b, -\gamma)
+ \hat{h}(\xi) - 2\pi i \xi_1 \gamma \hat{\eta}(\xi), \tag{3.64}
\]
which verifies the overdetermined compatibility condition (3.25).

Now we are ready to prove that \( \Upsilon_{\gamma, \kappa} : X^s \to \mathcal{Y}^s \) is an isomorphism when \( \sigma > 0 \) and \( n \geq 2 \), and when \( \sigma = 0 \) and \( n = 2 \).

**Theorem 3.16 (Existence and uniqueness of solutions to (3.28)).** Suppose \( \gamma \in \mathbb{R} \setminus \{0\}, s \geq 0 \).

1. (Isomorphism in the case with surface tension) If \( \sigma > 0 \) and \( n \geq 2 \), then the bounded linear map \( \Upsilon_{\gamma, \kappa} : X^s \to \mathcal{Y}^s \) defined in (3.44) is an isomorphism.
2. (Isomorphism in the case without surface tension) If \( \sigma = 0 \) and \( n = 2 \), then the bounded linear map \( \Upsilon_{\gamma, 0} : X^s \to \mathcal{Y}^s \) defined in (3.44) is an isomorphism.

**Proof.** To prove the first item, by Proposition 3.13, it suffices to show that \( \Upsilon_{\gamma, \kappa} \) is surjective. Suppose \( (f, g, h, k) \in \mathcal{Y}^s \), and define the free surface function \( \eta \in X^{s+5/2}(\Sigma; \mathbb{R}) \) by the construction in Lemma 3.11. By Theorem 3.9, there exists \( (u, p) \in \mathcal{H}^{s+2}(\Omega; \mathbb{R}^n) \times \mathcal{H}^{s+1}(\Omega; \mathbb{R}) \) such that \( \Psi_{\gamma}(u, p) = (\text{div } S(p, u) - \gamma \partial_1 u, \text{div } u, u_n|_{\Sigma_b}, S(p, u)e_n|_{\Sigma_b}) = (f - (\nabla \eta, 0), g, h - \gamma \partial_1 \eta, k - \sigma \Delta \eta e_n) \). Therefore, we find that \( \Upsilon_{\gamma, \kappa, \sigma}(u, p, \eta) = (f, g, h, k) \). This shows that \( \Upsilon_{\gamma, \kappa} \) is surjective, and it follows that \( \Upsilon_{\gamma, \kappa, \sigma} \) is an isomorphism.

To prove the second item we follow the same argument as above, using Lemma 3.15 in place of Lemma 3.11 and \( \Upsilon_{\gamma, 0} \) in place of \( \Upsilon_{\gamma, \kappa} \).

3.6. **Parameter regime for the linear isomorphism associated to (1.24).** In this subsection we consider the linearization of the flattened system (1.24) around the trivial solution \( u = 0, p = 0, \eta = 0 \), which is given by the map \( X^s \ni (u, p, \eta) \to L_{\kappa, \sigma} \in \mathcal{Y}^s \) defined in (1.42). Our goal is to show that there is a parameter regime in terms of \( \kappa \) for which \( L_{\kappa, \sigma} \) remains an isomorphism between \( X^s \) and \( \mathcal{Y}^s \). To achieve this we first define \( L_{\kappa, \sigma} : X^s \to \mathcal{Y}^s \) via
\[
L_{\kappa, \sigma}(u, p, \eta) = \Upsilon_{\gamma, \kappa, \sigma}(u, p, \eta) = (\text{div } S(p, u) - \gamma \partial_1 u + (\nabla \eta, 0), \text{div } u, u_n|_{\Sigma_b} + \gamma \partial_1 \eta, S(p, u)e_n|_{\Sigma_b} + \sigma \Delta \eta e_n), \tag{3.65}
\]
and
\[
M_{\kappa}(u, p, \eta) = (-\gamma \kappa x_n \partial_1 \eta e_1 + \kappa s_0(x_n) \partial_1 u + \kappa^2 s_0(x_n) u_n e_1 + \kappa^2 x_n s_0(x_n) \partial_1 \eta e_1 - \kappa x_n \Delta \eta e_1 - \kappa(x_n \nabla \partial_1 \eta, \partial_1 \eta),
\]
\[
\kappa x_n \partial_1 \eta, -\frac{k\kappa^2}{2} \partial_1 \eta, 0), \tag{3.66}
\]
so that \( L_{\kappa, \sigma} = L_{\kappa, \sigma}(I + L_{\kappa, \sigma}^{-1} M_{\kappa}) \), where \( L_{\kappa, \sigma} \) is defined via (1.42) and \( s_0 \) is defined via (1.6). Our goal in this subsection is to show that \( \|M_{\kappa}\| < \frac{1}{\|L_{\kappa, \sigma}\|} \) for sufficiently small \( \kappa \).

We first establish a preliminary lemma. While the following result is not at all surprising and follows from routine applications of Tonelli’s theorem and Hölder’s inequality, we record the proof below to show that the universal constants in the estimates are purely combinatorial and do not depend on the physical parameters.

**Lemma 3.17.** Let \( \mathbb{Z} \ni n \geq 2 \) and suppose \( f \in C^\infty([0, b]; \mathbb{R}) \). The following hold.
Proof. To prove the first item, we note that by Tonelli’s theorem and Hölder’s inequality,
\[
\|\partial^j f \partial^i g\|_{L^2} = \left( \int_\Sigma \int_0^b |\partial^j f(x,n)|^2 |\partial^i g(x')|^2 \, dx \, dx' \right)^{\frac{1}{2}} \leq \|\partial^j f\|_{L^\infty(0,b)} \left( \int_\Sigma \int_0^b |\partial^i g(x')|^2 \, dx \, dx' \right)^{\frac{1}{2}} \\
\leq \|\partial^j f\|_{L^\infty(0,b)} \|\partial^i g\|_{L^2(\Sigma)} \leq \|\partial^j f\|_{L^\infty} \|g\|_{H^k}. \tag{3.67}
\]

Then by the Leibniz formula, we have
\[
\|f g\|_{H^k} \leq \sum_{0 \leq |j| \leq k, \, m \in \mathbb{N}^n} \binom{j}{m} \|\partial^j f\|_{L^\infty} \|\partial^{j-m} g\|_{L^2} \leq C \max_{0 \leq j \leq k} \|\partial^j f\|_{L^\infty} \|g\|_{H^k}. \tag{3.68}
\]

where \(C = C(k) > 0\) is a combinatorial constant depending only on \(k\). To prove the second item, we note that by Tonelli’s theorem and Hölder’s inequality again,
\[
\|\partial^j f \partial^i g\|_{L^2} = \left( \int_\Sigma \int_0^b |\partial^j f(x,n)|^2 |\partial^i g(x')|^2 \, dx \, dx' \right)^{\frac{1}{2}} = \|\partial^j f\|_{L^2} \|\partial^i g\|_{L^2} \leq \|\partial^j f\|_{L^2} \|g\|_{H^k}. \tag{3.69}
\]

Then by the Leibniz formula,
\[
\|f g\|_{H^k} \leq \sum_{0 \leq |j| \leq k, \, m \in \mathbb{N}^n} \binom{j}{m} \|\partial^j f\|_{L^2} \|\partial^{j-m} g\|_{L^2} \leq C \max_{0 \leq j \leq k} \|\partial^j f\|_{L^2} \|g\|_{H^k}. \tag{3.70}
\]

This gives us an immediate corollary that will prove useful in the next section.

**Corollary 3.18** (Multiplication with smooth functions in the vertical variable). Let \(\Sigma\) be defined as in (1.1) and suppose \(R \ni s \geq 0\). If \(\varphi \in C^\infty([0,b];\mathbb{R})\), then the map \(T_{\varphi} : H^s(\Sigma;\mathbb{R}) \to H^s(\Omega;\mathbb{R})\) defined via
\[
T_{\varphi} f(x) = \varphi(x_n) f(x') \tag{3.71}
\]
is well-defined and bounded.

**Proof.** This follows directly from the second item of Lemma 3.17 when \(s\) is an integer. By interpolation, the second item also holds when \(s\) is real-valued. \(\square\)

Now we are ready to prove the main result of this subsection.

**Theorem 3.19.** Let \(R \ni s \geq 0, \kappa \in \mathbb{R}\) and consider the linear map \(M_\kappa : X^s \to Y^s\) as defined in (3.66). Then the following hold.

1. There exists a constant \(C > 0\) depending on \(s, n, \) and in the periodic cases on \(L_i\), such that
   \[
   \|M_\kappa\|_{L(X^s,Y^s)} \leq C \left( \text{max}\{b^2, b^3\} |\kappa|^2 + \left( |\gamma| \text{max}\{b^{1/2}, b^{3/2}\} + \text{max}\{1, b^2\} |\kappa| \right) \right). \tag{3.72}
   \]

2. There exists \(\kappa_0 > 0\) depending on \(\gamma, s, b, n, \) and in the toroidal cases on \(L_i\), such that for all \(\kappa \in (-\kappa_0, \kappa_0)\), the map \(L_{\kappa, \sigma} : X^s \to Y^s\) as defined in (3.65) is an isomorphism.
Proof. To prove the first item, we first suppose that \((u, p, \eta) \in \mathcal{X}^s\) for some integer \(s \geq 0\) and consider \((f, g, h, k) = M_{\kappa}(u, p, \eta)\), where \(M_{\kappa}\) is defined in (3.60). Then by the fifth item of Theorem 2.5, the third item of Theorem 2.6, and Lemma 3.17 we have \((f, g, h, k) = M_{\kappa}(u, p, \eta) \in \mathcal{Y}^s\). Recall that \(s_0(x_n) = bx_n - \frac{c^2}{2}\) and note that \(s_0(x_n) = 0\) for \(i \geq 3\) and \((x_n s_0(x_n))^{(i)} = 0\) for \(i \geq 4\). Therefore, by Lemma 3.17 and the fifth item of Theorem 2.5 we have

\[
\|f\|_{H^s} \lesssim_s \|\|\| \kappa \| \max_{0 \leq i < 1} \|x_n^{(i)}\|_{L^2} \|\partial_1 \eta\|_{H^s} + \|\|\| \kappa \| \max_{0 \leq i < 1} \|x_n^{(i)}\|_{L^2} \|\partial_1 \eta\|_{H^{s+1}} + \|\|\| \kappa \| \max_{0 \leq i < 1} \|x_n^{(i)}\|_{L^2} \|\partial_1 \eta\|_{H^{s+2}}
\]

\[\lesssim_{s, n} \left( \max\{b^+, b^-\} \|\kappa\|^2 + \left( \|\kappa\| \max\{b^{1/2}, b^{3/2}\} + \max\{1, b^2\} \right) \|\kappa\| \right) \|(u, p, \eta)\|_{\mathcal{X}^s}.
\]

We also have

\[
\|g\|_{H^{s+1}} \lesssim_s \|\|\| \kappa \| \max\{b^+, b^-\} \|\partial_1 \eta\|_{H^{s+1}} \lesssim_{s, n} \|\|\| \kappa \| \max\{b^+, b^-\} \|\eta\|_{\mathcal{X}^{s+2}}
\]

and \(\|h\|_{H^{s+\frac{1}{2}}} \lesssim_s \|\|\| \kappa \| \|\eta\|_{\mathcal{X}^{s+\frac{1}{2}}}\). Therefore,

\[
\|M_{\kappa}(u, p, \eta)\|_{\mathcal{Y}^s} \lesssim_{s, n} \left( \max\{b^+, b^-\} \|\kappa\|^2 + \left( \|\kappa\| \max\{b^{1/2}, b^{3/2}\} + \max\{1, b^2\} \right) \|\kappa\| \right) \|(u, p, \eta)\|_{\mathcal{X}^s}
\]

\[=: B(\gamma, b, \kappa) \|(u, p, \eta)\|_{\mathcal{X}^s} \tag{3.79}
\]

which implies that \(\|M_{\kappa}\|_{\mathcal{L}(\mathcal{X}^s, \mathcal{Y}^s)} \lesssim_{s, n} B(\gamma, b, \kappa)\). We note that by the fifth item of Theorem 2.5, in the toroidal cases the universal constants in the preceding equations also depend on \(L_i\).

For any real valued \(s \geq 0\), we may find an integer \(m \geq 0\) and some \(\theta \in [0, 1]\) such that \(s = \theta m + (1 - \theta)(m + 1)\). Then by interpolation, there exists some constant \(C > 0\) depending on \(s, n\), and in the periodic cases on \(L_i\) such that

\[
\|M_{\kappa}\|_{\mathcal{L}(\mathcal{X}^s, \mathcal{Y}^m)} \lesssim_{s, n} \|M_{\kappa}\|_{\mathcal{L}(\mathcal{X}^{m+1}, \mathcal{Y}^{m+1})} \leq CB(\gamma, b, \kappa).
\]

This proves the first item. To prove the second item, we note that

\[
B(\gamma, b, \kappa) = \max\{b^+, b^-\} \|\kappa\|^2 + \left( \|\kappa\| \max\{b^{1/2}, b^{3/2}\} + \max\{1, b^2\} \right) \|\kappa\| = \begin{cases} b^{7/2} \|\kappa\| + (\|\kappa\| b^{5/2} + b^2) \|\kappa\|, & b > 1 \\ b^{1/2} \|\kappa\|^2 + (|\kappa| b^{1/2} + 1) \|\kappa\|, & 0 < b \leq 1 \end{cases}.
\]

From (3.81), we may infer that for \(c = C^{-1} \|L_0, \sigma\|^{-1} > 0\), where \(C > 0\) is the constant appearing in (3.80), there exists \(\kappa_0\) depending on \(\gamma, s, b, n\) and in the toroidal cases on \(L_i\) for which \(B(\gamma, b, \kappa) < c\) if \(\kappa \in (\kappa_0, \kappa_0)\).

Then by (3.80), if \(\kappa \in (-\kappa_0, \kappa_0)\) then \(\|M_{\kappa}\|_{\mathcal{L}(\mathcal{X}^s, \mathcal{Y}^m)} < \|L_0, \sigma\|^{-1} \|\mathcal{L}(\mathcal{Y}^m, \mathcal{X}^s)\|^{-1}\), thus the linear map \(L_{0, \sigma}^{-1} M_{\kappa} : \mathcal{X}^s \to \mathcal{X}^s\) as defined in (3.65) satisfies \(\|L_{0, \sigma}^{-1} M_{\kappa}\| \leq \|L_0, \sigma\|^{-1} \|M_{\kappa}\| < 1\). Hence, \(I + L_{0, \sigma}^{-1} M_{\kappa} : \mathcal{X}^s \to \mathcal{X}^s\) is an isomorphism, and we conclude that \(L_{\kappa, \sigma} = L_{0, \sigma}(I + L_{0, \sigma}^{-1} M_{\kappa})\) is an isomorphism from \(\mathcal{X}^s\) to \(\mathcal{Y}^s\).

\[\square\]

4. Nonlinear analysis

4.1. Preliminaries. We first record a set of results on various maps defined in terms of \(\eta\), including the flattening map \(\mathfrak{F}_\eta\), that we will use in the subsequent analysis.

**Theorem 4.1.** Let \(\Sigma\) be defined as in (1.1), \(N \ni n \geq 2\), \(N \ni k > \frac{q}{2}\), and \(V\) be a real finite dimensional inner product space.

1. Let \(\zeta \in C_b^{0,1}(\Sigma; \mathbb{R})\) such that \(\inf \zeta > 0\). Then there exists \(r_1 > 0\) depending on \(n, b, s, \zeta,\) and in the toroidal cases on \(L_i\), such that the maps \(\Gamma_1, \Gamma_2 : B_{\mathcal{X}^s(0, r_1)} \times H^s(\Omega; \mathbb{R}) \to H^s(\Omega; \mathbb{R})\) given by

\[
\Gamma_1(f, g) = \frac{g}{b + f}, \quad \Gamma_2(f, g) = \frac{gf}{b + f}
\]

are well-defined and smooth. There also exists a constant \(r_2 > 0\) depending on \(n, b, s, \zeta,\) and in the toroidal cases on \(L_i\), such that the map \(\Gamma : B_{\mathcal{X}^s(0, r_2)} \to H^s(\Omega; \mathbb{R}^n)\) given by

\[
\Gamma(f) = \frac{f}{\sqrt{1 + |f|^2}}
\]

is well-defined and smooth.
Let \( \eta \in X^{k+\frac{3}{2}}(\Sigma; \mathbb{R}) \) be such that \( \|\eta\|_{C^b_{p}} \leq \frac{b}{2} \). Define \( \mathcal{G}_\eta : \Omega_{b} \to \Omega_{b} \) via

\[
\mathcal{G}_\eta(x) = \left( x', \frac{x_n b}{b + \eta(x')} \right).
\]

Then the following hold.

(a) \( \mathcal{G}_\eta \in C^r(\Omega_{b}; \Omega_{b}) \) is a diffeomorphism for \( r = 3 + \left[ k - \frac{n}{2} \right] \), with its inverse being \( \tilde{\mathcal{G}}_\eta \in C^r(\Omega; \Omega_{b}) \).

(b) If \( 0 \leq s \leq k + 2 \) and \( F \in H^s(\Omega_{b}; V) \), then \( F \circ \mathcal{G}_\eta \in H^s(\Omega_{b}; V) \). Moreover, there exists a constant \( c > 0 \) depending on \( n, s, k, \|\eta\|_{X^{k+2}} \), and in the toroidal cases on \( L_i \), such that

\[
\|F \circ \mathcal{G}_\eta\|_{H^s(\Omega_{b}; V)} \leq c \|F\|_{H^s(\Omega_{b}; V)},
\]

and the map \( r \mapsto c(n, s, k, r) \) is non-decreasing.

(c) If \( 0 \leq s \leq k + 2 \) and \( F \in H^s(\Omega_{b}; V) \), then \( F \circ \tilde{\mathcal{G}}_\eta \in H^s(\Omega_{b}; V) \). Moreover, there exists a constant \( c > 0 \) depending on \( n, s, k, \|\eta\|_{X^{k+2}} \), and in the toroidal cases on \( L_i \), such that

\[
\|F \circ \tilde{\mathcal{G}}_\eta\|_{H^s(\Omega_{b}; V)} \leq c \|F\|_{H^s(\Omega_{b}; V)},
\]

and the map \( r \mapsto c(n, s, k, r) \) is non-decreasing.

(3) (\( \omega \)-lemma for compositions) For \( \eta \in X^{k+\frac{3}{2}}(\Sigma; \mathbb{R}) \) we define the flattening map \( \tilde{\mathcal{G}}_\eta \) as in (1.17). Then there exists some \( 0 < \delta_0 < 1 \) such that the following hold:

(a) The map \( \Lambda : H^{k+1}(\Sigma; \mathbb{R}; \mathbb{R}^n) \times B_{X^{k+\frac{3}{2}}}(0, \delta_0) \to H^k(\Omega_{b}; \mathbb{R}^d) \) given by

\[
\Lambda(f, \eta) = f \circ \tilde{\mathcal{G}}_\eta
\]

is well-defined and \( C^1 \), with \( D\Lambda(f, \eta)(g, \zeta) = \frac{\partial}{\partial \eta} f \circ \tilde{\mathcal{G}}_\eta \zeta + g \circ \tilde{\mathcal{G}}_\eta \).

(b) The map \( \mathcal{G}_b : H^{k+2}(\Sigma; \mathbb{R}; \mathbb{R}^n) \times B_{X^{k+\frac{3}{2}}}(0, \delta_0) \to H^{k+\frac{3}{2}}(\Sigma; \mathbb{R}^d) \) defined via

\[
\mathcal{G}_b(f, \eta) = f \circ \tilde{\mathcal{G}}_\eta|_{\Sigma_b}
\]

is well-defined and \( C^1 \), with \( D\mathcal{G}_b(f, \eta)(g, \zeta) = \frac{\partial}{\partial \eta} f \circ \tilde{\mathcal{G}}_\eta \zeta + g \circ \tilde{\mathcal{G}}_\eta \)|\(\Sigma_b\).

(c) The map \( \mathcal{T}_b : H^{k+2}(\Sigma; \mathbb{R}; \mathbb{R}^n) \times B_{X^{k+\frac{3}{2}}}(0, \delta_0) \to H^{k+1/2}(\Sigma; \mathbb{R}^d) \) defined via

\[
\mathcal{T}_b(f, \eta) = (f \circ \tilde{\mathcal{G}}_\eta|_{\Sigma_b})|_{\mathcal{N}},
\]

where \( \mathcal{N} \) is defined via (1.13), is well-defined and \( C^1 \).

Proof. The three items in the case \( \Sigma = \mathbb{R}^d \) for \( d \geq 1 \) follow from Theorems 5.16, 5.17, A.12, Corollary 5.21, and Proposition 7.4 in [18]; the proofs therein can be adopted with minimal modification to handle to cases when \( \Sigma \neq \mathbb{R}^d \).

Now we can synthesize the aforementioned results to show that all the nonlinear maps appearing in (1.24) are well-defined and smooth.

**Theorem 4.2.** Let \( \mathbb{R} \ni s > \frac{n}{2} \) and \( \sigma > 0 \). For \( \delta > 0 \), define the open set of \( \mathcal{X}^s \)

\[
U_\delta^s = \{(u, p, \eta) \in \mathcal{X}^s : \|\eta\|_{X^{s+\frac{3}{2}}} < \delta \}
\]

For \( \gamma \in \mathbb{R}, T \in H^{s+\frac{3}{2}}(\Sigma; \mathbb{R}^{nxn}), (u, p, \eta) \in U_\delta^s, U_1^s = W_1 \circ \tilde{\mathcal{G}}_\eta, U_2^s = W_2 \circ \tilde{\mathcal{G}}_\eta \) with \( W_1, W_2 \) as defined in (1.7), (1.14), \( \tilde{\mathcal{G}}_\eta \) as defined in (1.17), and \( \mathcal{F}, \mathcal{A}, \mathcal{N}, \mathcal{H} \) as defined in (1.20), (1.21), (1.13) and (1.2), we define \( f = \text{div}_A S_A(p, u) - \gamma e_1 \cdot \nabla_A u - \gamma (x_n + \eta \frac{\partial}{\partial x_n}) \partial_1 \eta e_1 + (u + U^1 + U^2) \cdot \nabla_A (u + U^1 + U^2) + (\nabla \eta, 0) - \gamma (x_n + \eta \frac{\partial}{\partial x_n}) \Delta \eta e_1 - \gamma ((x_n + \eta \frac{\partial}{\partial x_n}) \nabla' \partial_1 \eta, \partial_1 \eta), g = \mathcal{F}[\text{div}_A u + \gamma (x_n + \eta \frac{\partial}{\partial x_n}) \partial_1 \eta(x')], h = u \cdot \mathcal{N} - (\gamma + s(\eta + b) + \sigma(\eta + b)\eta) \partial_1 \eta, \) and \( k = S_A(p, u) \mathcal{N} - [\sigma \mathcal{H}(\eta) I + T + \kappa(\eta + b)](e_1 \otimes (\nabla \eta, 0) + (\nabla \eta, 0) \otimes e_1)] \mathcal{N} \).

Then \( (f, g, h, k) \in \mathcal{Y}^s \), and the map

\[
\mathbb{R} \times \mathbb{R} \times H^{s+\frac{3}{2}}(\Sigma; \mathbb{R}^{nxn}) \times U_\delta^s \ni (\gamma, T, u, p, \eta) \mapsto (f, g, h, k) \in \mathcal{Y}^s
\]

is smooth.

Proof. By the first item in Theorem 2.6 and standard Sobolev embedding, there exists a constant \( \delta_0 > 0 \) depending on \( n, s, b, \) and in the toroidal cases on \( L_i \), such that if \( s > n/2 \) and \( \eta \in X^{s+5/2}(\Sigma; \mathbb{R}) \) with \( \|\eta\|_{X^{s+5/2}} < \delta_0 \), then \( \|\eta\|_{C^b_{p}} \leq b/2 \). We define \( \delta = \min\{\delta_0, r_1, r_2, c_1\} \), where \( r_1, r_2 \) are the radii from the first item of Theorem 4.1, and \( c_1 \) is the embedding constant from (2.24). By Theorem 7.3 in [18] and standard results from the theory of Sobolev spaces, the map

\[
\mathbb{R} \times \mathbb{R} \times H^{s+\frac{3}{2}}(\Sigma; \mathbb{R}^{nxn}) \times U_\delta^s \ni (\gamma, T, u, p, \eta) \mapsto (\text{div}_A S_A(p, u) + (u - \gamma e_1) \cdot \nabla_A u + u \cdot \nabla_A u, \mathcal{F}[\text{div}_A u, u \cdot \mathcal{N} + \gamma \partial_1 \eta, S_A(p, u) \mathcal{N} - (\sigma \mathcal{H}(\eta) I + T) \mathcal{N}] \in \mathcal{Y}^s
\]
is well-defined and smooth. Similarly, the maps functions in of Theorem 2.6, and Corollary 3.18 the map we have shown that \((\cdot, x_n) e_n = x_n + \eta(x') \frac{\partial \eta}{\partial x} \) is well-defined and smooth. By the fifth item of Theorem 2.5, the second item of Theorem 2.6, and Corollary 3.18 the map

\[
\mathbb{R} \times \mathbb{R} \times X^{s+5/2}(\Sigma; \mathbb{R}) \ni (\gamma, \eta) \mapsto -\kappa \mathfrak{H} \eta(x', x_n) e_n \partial_1 \eta e_1 + (\nabla' \eta, 0)
\]

is well-defined and smooth. By (1.7), (1.14), (1.17), and (1.21), we have

\[
u \cdot \nabla \mathcal{A}(U^1 + U^2) = \sum_{j=1}^{n-1} \left( \sum_{k=1}^{n} \frac{(U_1^1 + U_1^2 - 1) \mathcal{A}_{nk} \partial_k (U_1^1 + U_1^2) + (U_1^1 + U_1^2) \partial_1 (U_1^1 + U_1^2)}{b} \right)
\]

where \( \mathcal{K} = 1 + J \). By (1.7), (1.14), and (1.17), \((U_1 + U_2)(x, \eta) = -\kappa (b \mathfrak{H} \eta(x) e_n - (\mathfrak{H} \eta(x) e_n)^3) \frac{\partial \eta}{\partial x} \) is well-defined and smooth. By the first item of Theorem 4.1, the map \( B_{X^s}(0, \delta) \times H^s(\Omega; \mathbb{R}) \ni (\eta, \psi) \mapsto \psi \mathcal{K} \in H^s(\Omega; \mathbb{R}) \) is well-defined and smooth for any \( \psi \in H^s(\Omega; \mathbb{R}) \). We also note that every non-trivial term in the components of (1.14) is either a product of functions in \( X^{s+5/2}(\Sigma; \mathbb{R}) \) and functions in \( H^s(\Omega; \mathbb{R}) \), or functions in \( X^{s+5/2}(\Sigma; \mathbb{R}) \) and derivatives of functions in \( X^{s+5/2}(\Sigma; \mathbb{R}) \). To summarize, by the observations made above, (1.15), the second item of Theorem 2.6, Corollary 3.18, and the first item of Theorem 4.1, the map

\[
\mathbb{R} \times X^s \ni (\kappa, u, p, \eta) \mapsto u \cdot \nabla \mathcal{A}(U^1 + U^2) + (U^1 + U^2) \cdot \nabla \mathcal{A}(u + U^1 + U^2) \in H^s(\Omega; \mathbb{R}^n)
\]

is well-defined and smooth. Similarly, the maps

\[
\mathbb{R} \times X^{s+5/2}(\Sigma; \mathbb{R}) \ni (\kappa, \eta) \mapsto \mathcal{J} \kappa(x_n + \eta(x') \frac{\partial \eta}{\partial x}) \partial_1 \eta(x') \in H^{s+3/2}(\Sigma; \mathbb{R})
\]

and

\[
\mathbb{R} \times X^{s+5/2}(\Sigma; \mathbb{R}) \ni (\kappa, \eta) \mapsto (s(\kappa + b) \partial_1 \eta - \kappa(\kappa + b) \eta \partial_1 \eta, \kappa(\kappa + b) \eta, e_1 \odot (\nabla' \eta, 0) + (\nabla' \eta, 0) \odot e_1) \mathcal{N} \in \mathcal{H}^3(\Sigma; \mathbb{R}^n)
\]

are also well-defined and smooth. Finally, it remains to show that

\[
-s(\kappa + b) \partial_1 \eta - \kappa(\kappa + b) \eta \partial_1 \eta - \mathcal{J}(\cdot, x_n) \kappa \left( x_n + \eta(\cdot \frac{\partial \eta}{\partial x}) \right) \partial_1 \eta(\cdot) \, dx_n \in H^1(\Sigma; \mathbb{R})
\]

By (1.20), we have

\[
-s(\kappa + b) \partial_1 \eta - \kappa(\kappa + b) \eta \partial_1 \eta - \mathcal{J}(\cdot, x_n) \kappa \left( x_n + \eta(\cdot \frac{\partial \eta}{\partial x}) \right) \partial_1 \eta(\cdot) \, dx_n = -\kappa(b^2 \partial_1 \eta + b \partial_1 \eta^2 + \frac{1}{3} \partial_1 \eta^3).
\]

Thus, by Theorem 1.1 and the third item of Theorem 2.6 we have \( b^2 \partial_1 \eta + b \partial_1 \eta^2 + \frac{1}{3} \partial_1 \eta^3 \in H^{-1}(\Sigma; \mathbb{R}) \). Therefore, we have shown that \( (f, g, h, k) \in \mathcal{Y}^s \), and the map defined by (4.10) is smooth.
4.2. Solvability of the flattened system (1.24). Now we are ready to construct solutions to (1.24) by using the implicit function theorem.

**Proof of Theorem 1.2.** We first consider the case with surface tension, $\sigma > 0$ and $n \geq 2$. Let $\delta$ be the minimum of the $\delta_1 > 0$ from Theorem 4.2 and $\delta_* > 0$ from the third item of Theorem 4.1. Consider the open subset $U_\delta^* \subseteq \mathcal{X}^*$ defined via

$$U_\delta^* = \{(u,p,\eta) \in \mathcal{X}^* \mid \|\eta\|_{\mathcal{X}^*} < \frac{\delta}{2}\},$$

(4.21)

Using Proposition 3.11 and standard Sobolev embedding, any open subset of $U_\delta^*$ containing $(0,0,0)$ satisfies the first assertion of the theorem. This proves the first item.

To prove the remaining items, we consider the Hilbert space

$$\mathcal{E}^* = \mathbb{R} \times \mathbb{R} \times H^{s+2}(\Sigma \times \mathbb{R}; \mathbb{R}^{n \times n}) \times H^{s+\frac{2}{3}}(\Sigma; \mathbb{R}^{n \times n}) \times H^{s+1}(\Sigma \times \mathbb{R}; \mathbb{R}^n) \times H^s(\Sigma; \mathbb{R}^n)$$

(4.22)

and the solution map $\Xi : \mathcal{E}^* \times U_\delta^* \rightarrow \mathcal{Y}^*$ to be the derivatives of $\Xi$ with respect to $\mathcal{E}^*$ and $U_\delta^*$, respectively. Note that by the second item of Theorem 4.1, we have $D_2\Xi(0,0) = 0$ and $D_2 \Phi(0,0) = 0$. Therefore, for any $\gamma \in \mathbb{R}$, $\Xi(\gamma,0,0,0,0,\eta,0) = (0,0,0,0,0,0)$ and the second item of Theorem 4.1, we have $D_2 \Xi(\gamma,0,0,0,0,0,0,0) =\Xi_0$, where $\Xi_0$ is defined in (3.65). By Theorem 3.19, for every $\gamma \neq 0$ there exists some $\kappa_0 > 0$ for which $D_2 \Xi(\gamma,0,0,0,0,0,0,0) =\Xi_0$, where $\Xi_0$ is a linear isomorphism for every $\kappa$ in $(-\kappa_0, \kappa_0)$. Thus, by the implicit function theorem, there exists open sets $\mathcal{U}(\gamma) \subseteq \mathcal{E}^*$ and $\mathcal{O}(\gamma) \subseteq U_\delta^*$ such that $(\gamma,0,0,0,0,0) \in \mathcal{U}(\gamma)$ and $(0,0,0) \in \mathcal{O}(\gamma)$, and a $C^1$ and Lipschitz map $\varpi_{\gamma} : \mathcal{U}(\gamma) \rightarrow \mathcal{O}(\gamma)$ such that $\mathcal{E}(\gamma,\kappa,T,T,f,\varpi_{\gamma}(\gamma,\kappa,T,T,f)) = (0,0,0,0,0,0)$ for all $(\gamma,\kappa,T,T,f) \in \mathcal{U}(\gamma)$. Moreover, $(u,p,\eta) = \varpi_{\gamma}(\gamma,\kappa,T,T,f)$ is the unique solution to $\mathcal{E}(\gamma,\kappa,T,T,f,\varpi_{\gamma}(\gamma,\kappa,T,T,f)) = (0,0,0,0,0,0)$ in $\mathcal{O}(\gamma)$.

Next, we define the open sets

$$\mathcal{U}^* = \bigcup_{\gamma \in \mathbb{R} \setminus \{0\}} \mathcal{U}(\gamma) \subseteq \mathcal{E}^* \text{ and } \mathcal{O}^* = \bigcup_{\gamma \in \mathbb{R} \setminus \{0\}} \mathcal{O}(\gamma) \subseteq U_\delta^*.$$

(4.24)

We note that by construction, $(\mathbb{R} \setminus \{0\}) \times \{0\} \times \{0\} \times \{0\} \times \{0\} \subseteq \mathcal{U}^*$. Furthermore, for every $(\gamma,\kappa,T,T,f) \in U_\delta^*$, there exists a $\gamma_* \in \mathbb{R} \setminus \{0\}$ for which $(\gamma,\kappa,T,T,f) \in \mathcal{U}(\gamma_*)$ and $(u,p,\eta) = \varpi_{\gamma}(\gamma,\kappa,T,T,f) \in \mathcal{O}(\gamma_*)$. By the observation above and the implicit function theorem, the map $\varpi : \mathcal{U}^* \rightarrow \mathcal{O}^*$ defined via $\varpi(\gamma,\kappa,T,T,f) = \varpi_{\gamma}(\gamma,\kappa,T,T,f)$, where $\gamma_* \in \mathbb{R} \setminus \{0\}$ is such that $(\gamma,\kappa,T,T,f) \in \mathcal{U}(\gamma_*)$, is well-defined, $C^1$, and locally Lipschitz. This proves the remaining items for $\sigma > 0$ and $n \geq 3$.

To prove the remaining items in the case without surface tension and $n = 2$, we argue along the same lines but use the second item of Theorem 3.16 instead of the first and use the isomorphism $L_{\kappa,0}$.

\[\square\]

4.3. Solvability of the unflattened system (1.16). We now examine the solvability of the system (1.16) in the original Eulerian coordinates.

**Proof of Theorem 1.3.** Consider the map $\varpi : \mathcal{U}^* \rightarrow \mathcal{O}^*$ constructed in Theorem 1.2. By Theorem 1.2, for every $(\gamma,\kappa,T,T,f) \in \mathcal{U}^*$ the tuple $O^* = (u,p,\eta) = \varpi(\gamma,\kappa,T,T,f)$ solves (1.24) classically. Since $(u,p,\eta) \in \mathcal{O}^*$, we have $\|\eta\|_{C^2_b} \leq b/2$, therefore by the second item in Theorem 4.1, the flattening map $\mathcal{F}$ and its inverse $\mathcal{F}^{-1}$ are both $C^{(s-n/2)+3}$ diffeomorphisms. Now we fix $(\gamma,\kappa,T,T,f) \in \mathcal{U}^*$ and set $(u,p,\eta) = \varpi(\gamma,\kappa,T,T,f)$, $v = u \circ \mathcal{F}^{-1}$, and $q = p \circ \mathcal{F}^{-1}$. Then by the second item of Theorem 4.1 and Sobolev embedding, we have $v \in H^{s+2}(\Omega_{0+T}; \mathbb{R}^n) \cap C^{(s-n/2)+3}(\Omega_{0+T}; \mathbb{R}^n)$ and $q \in H^{s+1}(\Omega_{0+T}; \mathbb{R}^n) \cap C_b^{(s-n/2)+1}(\Omega_{0+T}; \mathbb{R}^n)$. Since $(\mathcal{F}^{-1}(x))' = x'$, we have $(f \circ \mathcal{F} + L_{0+T} f) \circ \mathcal{F}^{-1} (x) = f(x) + L_{0+T} f(x)$ and $(T \circ \mathcal{F} + S_{0+T} T) \circ \mathcal{F}^{-1} (x) = T_{\Sigma_{0+T}} (x) + S_{0+T} T(x)$ for all $x \in \Omega_{0+T}$. Thus, if $(u,p,\eta)$ is a solution tuple to (1.24) then $(v,q,\eta)$ is a solution tuple to (1.16). The last item follows from the fact that $\varpi$ is locally Lipschitz.\[\square\]
A.1. Permutations of Cartesian products. Here we collect a number of simple tools and bits of notation related to Cartesian products, especially products of groups.

**Definition A.1.** Suppose that \( d \geq 1 \) and that for \( 1 \leq i \leq d \) we have a commutative monoid \( X_i \) with identity element \( 0_i \in X_i \). Further suppose that \( E \subseteq \{1, \ldots, d\} \) and \( x \in X \). We define \( x_E \in X \) via

\[
(x_E)_j = \begin{cases} x_j & \text{if } j \in E \\ 0_j & \text{if } j \notin E. \end{cases}
\]

(A.1)

The utility of this notation is seen in the formula \( x = x_E + x_{E^c} \), valid for any set \( E \subseteq \{1, \ldots, d\} \).

We will often use this notation when \( X = \Gamma \) or \( X = \hat{\Gamma} \), where \( \Gamma \) is of the form (1.43). In this context, we introduce further notation.

**Definition A.2.** Let \( \Gamma \) be as in (1.43).

1. We define the sets \( R_\Gamma, R_\hat{\Gamma}, T_\Gamma, T_{\hat{\Gamma}} \subseteq \{1, \ldots, d\} \) via

\[
R_\Gamma = R_\hat{\Gamma} = \{i \in \{1, \ldots, d\} : \Gamma_i = \mathbb{R}\} =: \{\tau_1, \ldots, \tau_{R_\Gamma}\}
\]

(A.2)

\[
T_\Gamma = T_{\hat{\Gamma}} = \{i \in \{1, \ldots, d\} : \Gamma_i = L_\mathbb{T}\} = \{1, \ldots, d\} \setminus R_\Gamma =: \{t_1, \ldots, t_{T_\Gamma}\},
\]

(A.3)

where we order these so that \( \tau_1 < \cdots < \tau_{R_\Gamma} \) and \( t_1 < \cdots < t_{T_\Gamma} \). This allows us to write \( \Gamma \ni x = x_{R_\Gamma} + x_{T_\Gamma} \) and \( \Gamma \ni \xi = \xi_{R_\Gamma} + \xi_{T_\Gamma} \).

2. We define the sets \( \Gamma_R, \Gamma_T, \hat{\Gamma}_R, \hat{\Gamma}_T \) via

\[
\Gamma_R = \{x \in \Gamma \mid x_{T_\Gamma} = 0\}, \quad \Gamma_T = \{x \in \Gamma \mid x_{R_\Gamma} = 0\}, \quad \hat{\Gamma}_R = \{\xi \in \hat{\Gamma} \mid \xi_{T_\hat{\Gamma}} = 0\}, \quad \hat{\Gamma}_T = \{\xi \in \hat{\Gamma} \mid \xi_{R_\hat{\Gamma}} = 0\}.
\]

(A.4)

It follows that we have the direct sum decompositions

\[
\Gamma = \Gamma_R \oplus \Gamma_T, \quad \hat{\Gamma} = \hat{\Gamma}_R \oplus \hat{\Gamma}_T.
\]

(A.5)

At numerous points in the paper it is convenient to reorder the factors appearing in a Cartesian product. We introduce notation for this and related ideas now.

**Definition A.3.** Suppose that \( X_1, \ldots, X_d \) are sets and consider the Cartesian product \( X = \prod_{i=1}^d X_i \). Let \( S_d \) denote the symmetric group of permutations of the indices \( \{1, \ldots, d\} \).

1. For \( g \in S_d \) we write \( gX \) for the reordered Cartesian product \( gX = \prod_{i=1}^d X_{g(i)} \).

2. Let \( g \in S_d \). We define the induced reordering map \( P_g : X \to gX \) via

\[
P_g(x) = P_g(x_1, \ldots, x_d) = (x_{g(1)}, \ldots, x_{g(d)}).
\]

(A.6)

This map is clearly a bijection. Sometimes it will be convenient to include the domain in the notation, in which case we write \( P_{g,X} \) in place of \( P_g \).

3. Given sets \( E \) and \( Y \), write \( F(E; Y) = \{f : E \to Y\} \). Let \( g \in S_d \) and \( Y \) be a set. We define the map \( \sigma_g : F(X; Y) \to F(gX; Y) \) via \( \sigma_g f = f \circ P_g^{-1} \). Since \( P_g \) is invertible, the map \( \sigma_g \) is as well, and \( \sigma_g : F(X; Y) \to F(X; Y) \) is given by \( \sigma_g^{-1} f = f \circ P_g \). As above, we sometimes indicate the domain by writing \( \sigma_{g,X} \) in place of \( \sigma_g \).

It will be convenient to record the following obvious result as a lemma.

**Lemma A.4.** Suppose that \( X_1, \ldots, X_d \) are sets and consider the Cartesian product \( X = \prod_{i=1}^d X_i \). Let \( g \in S_d \) and let \( P_g \) be the induced reordering map defined in the second item in Definition A.3. Then the following hold.

1. If each \( X_i \) is a group, then \( P_g \) is a group isomorphism.

2. If each \( X_i \) is a topological space and \( X \) is endowed with the product topology, then \( P_g \) is a homeomorphism.

3. If each \( X_i \) is a metric space and \( X \) is endowed with a product metric, then \( P_g \) is an isometry.

4. If each \( X_i \) is a \( \sigma \)-finite measure space and \( X \) is endowed with the product measure, then \( P_g \) is a measure preserving map.

5. If each \( X_i \) is a smooth manifold and \( X \) is endowed with the smooth product manifold structure, then \( P_g \) is a smooth diffeomorphism.

Next we introduce the notion of the canonical reordering of \( \Gamma \) and \( \hat{\Gamma} \), where \( \Gamma \) is as in (1.43).

**Definition A.5.** Let \( \Gamma \) be as in (1.43), and let \( R_\Gamma, R_\hat{\Gamma}, T_\Gamma, T_{\hat{\Gamma}} \subseteq \{1, \ldots, d\} \) be as in Definition A.2.
Lemma A.6. Let \( \Gamma \) be defined as in (1.43) and assume \( 1 \in R_F \). We define the surjective map \( \pi_{\hat{\Gamma}} : \hat{\Gamma} \to \mathbb{R}^{[R_F]} \) via
\[
\pi_{\hat{\Gamma}} = \pi_{\hat{\Gamma}} \circ \mathcal{P}_{\hat{\Gamma}}.
\]

The following hold.

1. The restriction of the map \( \pi_{\hat{\Gamma}} \) on \( \hat{\Gamma}_{*,R} \) is an isometric measure-preserving group isomorphism between \( \hat{\Gamma}_{*,R} \) and \( \mathbb{R}^{[R_F]} \), where \( \hat{\Gamma}_{*,R} \) is the first summand in the direct sum decomposition of \( \hat{\Gamma}_{*} \) in (A.9).

2. As a consequence of the first item and Lemma A.4, the restriction of the map \( \pi_{\hat{\Gamma}} \) on \( \hat{\Gamma}_R \) is an isometric measure-preserving group isomorphism between \( \hat{\Gamma}_R \) and \( \mathbb{R}^{[R_F]} \), with the map \( \mathcal{P}_{\hat{\Gamma}}^{-1} \circ l_{\hat{\Gamma}} \) as its inverse.

3. The restriction of \( \pi_{\hat{\Gamma}} \) on \( B(0,r) \subseteq \hat{\Gamma}_R \) is also an isometric measure-preserving group isomorphism between \( B_F(0,r) \) to \( B_{\mathbb{R}^{[R_F]}}(0,r) \), with the restriction of the map \( \mathcal{P}_{\hat{\Gamma}}^{-1} \circ l_{\hat{\Gamma}} \) on \( B_{\mathbb{R}^{[R_F]}}(0,r) \) as its inverse.

In other words, we have \( \xi \in B_F(0,r) \Rightarrow \pi_{\hat{\Gamma}} \xi \in B_{\mathbb{R}^{[R_F]}}(0,r) \) and \( \omega \in B_{\mathbb{R}^{[R_F]}}(0,r) \Rightarrow \mathcal{P}_{\hat{\Gamma}}^{-1} l_{\hat{\Gamma}} \omega \in B_F(0,r) \), and the maps \( \pi_{\hat{\Gamma}}, \mathcal{P}_{\hat{\Gamma}}^{-1} l_{\hat{\Gamma}} \) are isometric measure-preserving group isomorphisms.
A.2. Tempered distributions and the Fourier transform. In this subsection we carefully define the class of tempered distributions on $\Gamma$ and its Pontryagin dual $\hat{\Gamma}$. First, we note that for any $\alpha \in \mathbb{N}^d$, we may follow the convention in (A.1) and write $\alpha = \alpha_{R_T} + \alpha_{T_T}$. We then define
\[
\mathbb{N}^d_{R_T} = \{ \alpha \in \mathbb{N}^d \mid \alpha_{R_T} = 0 \}, \quad \mathbb{N}^d_{T_T} = \{ \alpha \in \mathbb{N}^d \mid \alpha_{T_T} = 0 \}.
\]

**Definition A.7.** Let $\Gamma$ be defined as in (1.43).

1. We define the Schwartz class on $\Gamma$ to be
\[
\mathcal{S}(\Gamma; C) = \{ f \in C^\infty(\Gamma; C) \mid |f|_{\Gamma,\alpha,\beta} = \sup_{x \in \Gamma} |x^\alpha D^\beta f(x)| < \infty \text{ for all } \alpha \in \mathbb{N}^d_{R_T}, \beta \in \mathbb{N}^d \},
\]
and on $\hat{\Gamma}$ to be
\[
\mathcal{S}(\hat{\Gamma}; C) = \{ f : \hat{\Gamma} \to C \mid f(\mathcal{P}^{-1}_\Gamma(\cdot, 0) + \xi_T) \in C^\infty(\mathbb{R}^{|\xi_T|}; C) \text{ for all } \xi_T \in \hat{\Gamma}_{T_T} \text{ and } |f|_{\hat{\Gamma},\alpha,\beta} = \sup_{\xi \in \Gamma} |\xi^\alpha D^\beta f(\xi)| < \infty \text{ for all } \alpha \in \mathbb{N}^d, \beta \in \mathbb{N}^d \}.
\]

We equip $\mathcal{S}(\Gamma; C)$ and $\mathcal{S}(\hat{\Gamma}; C)$ with the Fréchet topology induced by the countable family of seminorms $\{ ||\cdot||_{\Gamma,\alpha,\beta} \}_{\alpha \in \mathbb{N}^d_{R_T}, \beta \in \mathbb{N}^d}$ and $\{ ||\cdot||_{\hat{\Gamma},\alpha,\beta} \}_{\alpha \in \mathbb{N}^d, \beta \in \mathbb{N}^d_{T_T}}$, by endowing $\mathcal{S}(\Gamma; C)$ and $\mathcal{S}(\hat{\Gamma}; C)$ with the metrics
\[
d_{\mathcal{S}(\Gamma)}(f, g) = \sum_{\alpha \in \mathbb{N}^d_{R_T}} \sum_{\beta \in \mathbb{N}^d} \frac{1}{2^{|\alpha|+|\beta|}} \frac{|f - g|_{\Gamma,\alpha,\beta}}{1 + |f - g|_{\Gamma,\alpha,\beta}}, \quad d_{\mathcal{S}(\hat{\Gamma})}(f, g) = \sum_{\alpha \in \mathbb{N}^d} \sum_{\beta \in \mathbb{N}^d_{T_T}} \frac{1}{2^{|\alpha|+|\beta|}} \frac{|f - g|_{\hat{\Gamma},\alpha,\beta}}{1 + |f - g|_{\hat{\Gamma},\alpha,\beta}}.
\]

2. We define the unitary Fourier and inverse Fourier transforms $\mathcal{F}_\Gamma^\pm : \mathcal{S}(\Gamma; C) \to \mathcal{S}(\hat{\Gamma}; C), \mathcal{F}_\Gamma^\pm : \mathcal{S}(\hat{\Gamma}; C) \to \mathcal{S}(\Gamma; C)$ on $\mathcal{S}(\Gamma; C)$ and $\mathcal{S}(\hat{\Gamma}; C)$ via
\[
\mathcal{F}_\Gamma^\pm \{ f \}(\xi) = \frac{1}{\sqrt{\prod_{T_T} L_i}} \int_{\Gamma} f(x) e^{\mp 2\pi i \xi \cdot x} \, dx = \mathcal{F}_\Gamma \{ f \}(\xi),
\]
\[
\mathcal{F}_\Gamma \{ f \}(x) = \frac{1}{\sqrt{\prod_{T_T} L_i}} \int_{\Gamma} f(\xi) e^{\mp 2\pi i \xi \cdot x} \, d\xi = \mathcal{F}_\Gamma^\mp \{ f \}(x).
\]

3. We define the class of tempered distributions $\mathcal{S}'(\Gamma; C)$ and $\mathcal{S}'(\hat{\Gamma}; C)$ to be the set of continuous linear functionals on $\mathcal{S}(\Gamma; C)$ and $\mathcal{S}(\hat{\Gamma}; C)$, respectively.

The next lemma shows that the restriction of $\sigma_{g,\Gamma}$ on $\mathcal{S}(\Gamma; F)$ is an isometric isomorphism between $\mathcal{S}(\Gamma; F)$ and $\mathcal{S}(g\Gamma; F)$, where $\sigma_{g,\Gamma}$ is defined in the third item of Definition A.3.

**Lemma A.8.** Let $\Gamma$ be defined as in (1.43), $g \in S_d$ be a permutation, and let $g\Gamma$ be as defined in the first item of Definition A.3. Then the map $\sigma_{g,\Gamma} : \mathcal{S}(\Gamma; F) \to \mathcal{S}(g\Gamma; F)$ defined in the third item of Definition A.3 is an isometric isomorphism, with its inverse given by the map $\sigma_{g,\Gamma}^{-1} : \mathcal{S}(g\Gamma; F) \to \mathcal{S}(\Gamma; F)$, also defined in the third item of Definition A.3.

**Proof.** For $g \in S_d$, let $\mathcal{P}_g, \Gamma : \Gamma \to g\Gamma, \mathcal{P}_g, \mathcal{N} : \mathbb{N}^d \to g\mathbb{N}^d$ be the induced reordering maps defined via (A.6). We note that for any $F \in \mathcal{S}(\Gamma; F), \alpha \in \mathbb{N}^d_{R_T},$ and $\beta \in \mathbb{N}^d$ we have $F(\mathcal{P}_g^{-1} \cdot) \in C^\infty(g\Gamma; F)$ with
\[
[F(\mathcal{P}_g^{-1} \cdot)]_{g\Gamma,\alpha,\beta} = \sup_{x \in g\Gamma} |x^\alpha D^\beta F(\mathcal{P}_g^{-1} x)| = \sup_{y \in \Gamma} |y^{g,\mathcal{N}}_{\cdot,\alpha} \partial_y^{g,\mathcal{N}}_{\cdot,\beta} F(y)| = [F]_{\Gamma,\alpha,\beta} < \infty,
\]
and for any $G \in \mathcal{S}(g\Gamma; F), \alpha \in \mathbb{N}^d_{R_T},$ and $\beta \in \mathbb{N}^d$ we have $G(\mathcal{P}_g^{-1} \cdot) \in C^\infty(\Gamma; F)$ with
\[
[G(\mathcal{P}_g^{-1} \cdot)]_{\Gamma,\alpha,\beta} = \sup_{x \in \Gamma} |x^\alpha D^\beta G(\mathcal{P}_g^{-1} x)| = \sup_{y \in \Gamma} |y^{g,\mathcal{N}}_{\cdot,\alpha} \partial_y^{g,\mathcal{N}}_{\cdot,\beta} G(y)| = [G]_{\Gamma,\alpha,\beta} < \infty.
\]
This shows that $\sigma_{g,\Gamma}$ and $\sigma_{g,\Gamma}^{-1}$ are well-defined. Since $\sigma_{g,\Gamma}$ is a bijection, and $\mathcal{P}_g, \mathcal{N}$ is a bijection between $\mathbb{N}^d_{R_T}$ and $\mathbb{N}^d_{R_{gT}},$ and between $\mathbb{N}^d$ and $g\mathbb{N}^d = \mathbb{N}^d,$ with (A.20) and (A.21) we can also conclude that
\[
d_{\mathcal{S}(\Gamma)}(F, G) = d_{\mathcal{S}(g\Gamma)}(F(\mathcal{P}_g^{-1} \cdot), G(\mathcal{P}_g^{-1} \cdot)).
\]
This shows that $\sigma_{g,\Gamma}$ is an isometric isomorphism. \qed
Following the exact same argument as above we can conclude that the restriction of \( \sigma_{g,f} \) on \( \mathcal{S}(\hat{\Gamma}; \mathcal{F}) \) is an isometric isomorphism between \( \mathcal{S}(\Gamma; \mathcal{F}) \) and \( \mathcal{S}(\hat{\Gamma}; \mathcal{F}) \).

Next, we record some important properties of the Schwartz classes \( \mathcal{S}(\Gamma; \mathcal{F}) \) and \( \mathcal{S}(\hat{\Gamma}; \mathcal{F}) \).

**Proposition A.9.** Let \( \Gamma \) be defined as in (1.43). Then the following hold.

1. We have the continuous inclusions \( \mathcal{S}(\Gamma; \mathcal{C}) \hookrightarrow L^p(\Gamma; \mathcal{C}) \) and \( \mathcal{S}(\hat{\Gamma}; \mathcal{C}) \hookrightarrow L^p(\hat{\Gamma}; \mathcal{C}) \) for any \( 1 \leq p \leq \infty \).
2. We have the continuous inclusion \( L^p(\Gamma; \mathcal{C}) \hookrightarrow \mathcal{S}(\Gamma; \mathcal{C}) \) and \( L^p(\hat{\Gamma}; \mathcal{C}) \hookrightarrow \mathcal{S}(\hat{\Gamma}; \mathcal{C}) \) for all \( 1 \leq p \leq \infty \).
3. For any \( f \in \mathcal{S}(\Gamma; \mathcal{C}) \), \( \alpha \in \mathbb{N}^d_{R_{\text{tr}}} \), and \( \beta \in \mathbb{N}^d \) the functions \( g_{\alpha,\beta} : \Gamma \to \mathbb{C} \) defined via \( g_{\alpha,\beta}(x) = x^\alpha \partial^\beta f(x) \) and \( h_{\alpha,\beta} = \partial^\alpha (x^\beta f(x)) \) belong to \( \mathcal{S}(\Gamma; \mathcal{C}) \). Furthermore, the maps \( f \mapsto g_{\alpha,\beta} \) and \( f \mapsto h_{\alpha,\beta} \) are continuous.

**Proof.** Define

\[
\tau = \begin{cases} 
1, & \Gamma = \mathbb{R}^d \\
\max\{1, 2 \max_{i \in \{1, \ldots, d\}} L_i\}, & \text{otherwise.}
\end{cases}
\] (A.24)

To prove the first item, we note that for any \( 1 \leq p < \infty \), \( f \in \mathcal{S}(\Gamma; \mathcal{C}) \), and \( \alpha, \beta \in \mathbb{N}^d_{R_{\text{tr}}} \) with \( |\alpha| = d+1, |\beta| = [(d+1)/p] + 1 \) we have

\[
\|f\|_{L^p(\Gamma)} \lesssim \left( \int_{B(0,\tau)} |f(x)|^p \, dx + \int_{B(0,\tau)^c} d_{\Gamma}(x,0)^\alpha |f(x)|^p d_{\Gamma}(x,0)^{-\alpha} \, dx \right)^{1/p} 
\lesssim \left( \|f\|_{L^\infty}^p + \left( \sup_{x \in B(0,\tau)^c} d_{\Gamma}(x,0)^\alpha |f(x)| \right) \int_{B(0,\tau)^c} d_{\Gamma}(x,0)^{-\alpha} \, dx \right)^{1/p} 
\lesssim \|f\|_{L^\infty} + \left( \sum_{\beta \in \mathbb{N}^d_{R_{\text{tr}}}, |\beta| = [(d+1)/p] + 1} [f]_{\Gamma, \beta, 0} \right)
\] (A.25)
and when \( p = \infty \) we have \( \|f\|_{L^\infty(\Gamma)} = [f]_{\Gamma, 0, 0} \). It follows then that \( \mathcal{S}(\Gamma; \mathcal{C}) \hookrightarrow L^p(\Gamma; \mathcal{C}) \) for all \( 1 \leq p \leq \infty \). A similar set of computations can be performed to prove \( \mathcal{S}(\hat{\Gamma}; \mathcal{C}) \hookrightarrow L^p(\hat{\Gamma}; \mathcal{C}) \) for all \( 1 \leq p \leq \infty \). For the second item, suppose that we have a sequence \( \{f_k\}_{k=1}^\infty \subset L^p(\Gamma; \mathcal{C}) \) for \( 1 \leq p \leq \infty \) such that \( f_k \to f \) in \( L^p(\Gamma; \mathcal{C}) \). We note that the canonical distribution \( T_{f_k} : \mathcal{S}(\Gamma; \mathcal{C}) \to \mathcal{C} \) associated to each \( f_k \) defined via

\[
\langle T_{f_k}, \varphi \rangle_{\mathcal{S}', \mathcal{S}} = \int_{\Gamma} f_k \varphi \, dx \text{ for all } \varphi \in \mathcal{S}(\Gamma; \mathcal{C})
\] (A.26)

is clearly linear. By the first item and Hölder’s inequality, if \( \{\varphi_m\}_{m=1}^\infty \subset \mathcal{S}(\Gamma; \mathcal{C}) \) and \( \varphi_m \to \varphi \) as \( m \to \infty \) in \( \mathcal{S}(\Gamma; \mathcal{C}) \), then for each \( k \) we have

\[
|\langle T_{f_k}, \varphi_m - \varphi \rangle_{\mathcal{S}', \mathcal{S}}| \leq \|f_k\|_{L^p} \|\varphi_m - \varphi\|_{L^q} \to 0.
\] (A.27)

This implies that \( T_{f_k} \in \mathcal{S}'(\Gamma; \mathcal{C}) \). Furthermore, since \( f_k \to f \) in \( L^p(\Gamma; \mathcal{C}) \), by the first item and Hölder’s inequality again, we have

\[
|\langle T_{f_k} - T_f, \varphi \rangle_{\mathcal{S}', \mathcal{S}}| \leq \|f_k - f\|_{L^p} \|\varphi\|_{L^q} \to 0 \text{ for all } \varphi \in \mathcal{S}(\Gamma; \mathcal{C}).
\] (A.28)
Since the map $L^p(\Gamma; \mathbb{C}) \ni f \mapsto T_f \in \mathcal{S}'(\Gamma; \mathbb{C})$ is also clearly injective, the inclusion map $i : L^p(\Gamma; \mathbb{C}) \to \mathcal{S}'(\Gamma; \mathbb{C})$ as is well-defined and continuous. This proves the second item for $\mathcal{S}'(\Gamma; \mathbb{C})$, and the same argument can be applied to prove the second item for $\mathcal{S}'(\hat{\Gamma}; \mathbb{C})$. The third item follows directly from the Leibniz rule.

For the remaining items, we note that by using the isometric isomorphisms $\sigma_{\Gamma} : \mathcal{S}(\Gamma; \mathbb{F}) \to \mathcal{S}(\hat{\Gamma}; \mathbb{F})$ and $\sigma_{\hat{\Gamma}} : \mathcal{S}'(\hat{\Gamma}; \mathbb{F}) \to \mathcal{S}'(\Gamma; \mathbb{F})$ induced by the canonical permutation $\varpi_{\gamma}$, it suffices to prove the remaining items for $\mathcal{S}(\hat{\Gamma}; \mathbb{F})$ and $\mathcal{S}'(\Gamma; \mathbb{F})$. Indeed, if the remaining items are true for $\mathcal{S}(\Gamma; \mathbb{F})$ and $\mathcal{S}'(\hat{\Gamma}; \mathbb{F})$, then it suffices to check that

\[ \mathcal{F}_{\Gamma}^\pm = \sigma_{\Gamma}^{-1} \circ \mathcal{F}_{\Gamma, r}^\pm \circ \sigma_{\Gamma}, \]  
\hspace{1cm} (A.29)

and the remaining items also follow for $\mathcal{S}(\Gamma; \mathbb{F})$ and $\mathcal{S}'(\hat{\Gamma}; \mathbb{F})$ as the following diagram commutes:

\[
\begin{array}{ccc}
\mathcal{S}(\Gamma; \mathbb{F}) & \xrightarrow{\pi_{\Gamma}} & \mathcal{S}(\hat{\Gamma}; \mathbb{F}) \\
\downarrow{\sigma_{\Gamma}} & & \downarrow{\sigma_{\Gamma}} \\
\mathcal{S}'(\Gamma; \mathbb{F}) & \xrightarrow{\pi_{\Gamma}^*} & \mathcal{S}'(\hat{\Gamma}; \mathbb{F})
\end{array}
\]  
\hspace{1cm} (A.30)

To verify (A.29), we note that for any $f \in \mathcal{S}(\Gamma; \mathbb{F})$, we have

\[
(\sigma_{\Gamma}^{-1} \mathcal{F}_{\Gamma, r}^\pm \sigma_{\Gamma} f)(\xi) = \sigma_{\Gamma}^{-1} \mathcal{F}_{\Gamma, r}^\pm [\mathcal{F}(\pi_{\Gamma}^{-1} f)](\xi) = \sigma_{\Gamma}^{-1} \frac{1}{\prod_{i \in \Gamma} L_i} \int_{\Gamma} f(x) e^{\mp 2\pi i \xi \cdot x} \, dx
\]

\[
= \sigma_{\Gamma}^{-1} \frac{1}{\prod_{i \in \Gamma} L_i} \int_{\Gamma} f(x) e^{\mp 2\pi i \xi \cdot x} \, dx = \sigma_{\Gamma}^{-1} \mathcal{F}_{\Gamma, r}^\pm [\mathcal{F}(f)](\xi) = \mathcal{F}_{\Gamma}^\pm [\mathcal{F}(f)](\xi). \]  
\hspace{1cm} (A.31)

A similar set of calculations can be performed to verify that $\sigma_{\Gamma}^{-1} \circ \mathcal{F}_{\Gamma, r}^\pm \circ \sigma_{\Gamma}$. Consequently, for the rest of the proof we may assume without loss of generality that $\Gamma = \mathbb{R}^{d_1} \times \prod_{i=1}^{d_2} L_i \mathbb{Z}$, with $d = \dim \Gamma = d_1 + d_2$.

To prove the fourth item, we first note that if $f \in \mathcal{S}(\Gamma; \mathbb{F})$, by the first and third items we have $\partial^\alpha ((-2\pi i)^{\gamma} f) \in \mathcal{S}(\hat{\Gamma}; \mathbb{C}) \subset L^1(\Gamma; \mathbb{C})$ for any $\alpha, \beta, \gamma \in \mathbb{N}^{d_1}$, $\beta \neq 0$. This shows that $\mathcal{F}_{\Gamma}^\pm (f, k) \in C^\infty(\mathbb{R}^{d_1}; \mathbb{C})$ for any $k \in \prod_{i=1}^{d_2} L_i^{-1} \mathbb{Z}$. Then for any $\alpha, \gamma \in \mathbb{N}^{d_1}, \beta \in \mathbb{N}^{d_2}$, we have

\[
(-2\pi i)^\alpha (-2\pi i)^{\beta} \partial^\alpha_x \mathcal{F}_{\Gamma}^\pm (f, \xi, k) = (-2\pi i)^{\alpha + \beta} \partial^\alpha_x \int_{\Gamma} f(x, y) e^{\mp 2\pi i \xi \cdot x} e^{\mp 2\pi i \beta y} \, dxdy
\]

\[
= \int_{\Gamma} f(x, y) e^{\mp 2\pi i \xi \cdot x} \partial_x^\alpha \partial_y^\beta f(x, y) e^{\mp 2\pi i \beta y} \, dxdy = (-1)^{|\alpha| + |\beta|} \int_{\Gamma} \partial_x^\alpha (e^{\mp 2\pi i \xi \cdot x}) \partial_y^\beta f(x, y) e^{\mp 2\pi i \beta y} \, dxdy. \]  
\hspace{1cm} (A.32)

By the first and third items, we find that the preceding inequality implies $\mathcal{F}_{\Gamma}^\pm \{f \} \in \mathcal{S}(\hat{\Gamma}; \mathbb{C})$ and if $\{ f_k \}_{\kappa=1}^\infty \subset \mathcal{S}(\Gamma; \mathbb{C})$ and $f_k \to f$ in $\mathcal{S}(\Gamma; \mathbb{C})$, then $\mathcal{F}_{\Gamma}^\pm \{f_k \} \to \mathcal{F}_{\Gamma}^\pm \{f \}$ in $\mathcal{S}(\Gamma; \mathbb{C})$. Therefore, we can conclude that if $f \in \mathcal{S}(\Gamma; \mathbb{C})$ then $\mathcal{F}_{\Gamma}^\pm \{f \} \in \mathcal{S}(\Gamma; \mathbb{C})$, and the map $f \mapsto \mathcal{F}_{\Gamma}^\pm \{f \}$ is continuous. Likewise, by a similar calculation we find that if $f \in \mathcal{S}(\Gamma; \mathbb{C})$ then $\mathcal{F}_{\Gamma}^\pm \{f \} \in \mathcal{S}(\Gamma; \mathbb{C})$, and the map $f \mapsto \mathcal{F}_{\Gamma}^\pm \{f \}$ is continuous. To prove the fifth item, we note that by the first item and Fubini’s theorem,

\[
\int_{\Gamma} \mathcal{F}_{\Gamma}^\pm \{f \}(\xi, k) g(\xi) \, d\xi \, dk = \int_{\Gamma} \left( \int_{\Gamma} f(x, y) e^{\mp 2\pi i \xi \cdot x} e^{\mp 2\pi i \beta y} \, dxdy \right) g(\xi) \, d\xi \, dk
\]

\[
= \int_{\Gamma} f(x, y) \left( \int_{\Gamma} g(\xi) e^{\mp 2\pi i \xi \cdot x} e^{\mp 2\pi i \beta y} \, d\xi \right) \, dk = \int_{\Gamma} f(x, y) \mathcal{F}_{\Gamma}^\pm \{g \}(x, y) \, dxdy. \]  
\hspace{1cm} (A.33)

To prove the sixth item, suppose $f \in \mathcal{S}(\Gamma; \mathbb{C})$ and for any fixed $t > 0$ and $(x, y) \in \Gamma$ we consider the function

\[ \phi^\pm(\xi, k) = e^{-\pi t^2 (|\xi|^2 + |k|^2)} e^{\mp 2\pi i \xi \cdot x} \frac{e^{\pm 2\pi i \beta y}}{\sqrt{\prod_{i=1}^{d_2} L_i}}. \]  
\hspace{1cm} (A.34)

We note that

\[ \mathcal{F}_{\Gamma}^\pm \{\phi^\pm \}(r, s) = \mathcal{F}_{\Gamma}^\pm \{g_t \}(r - x, s - y), \]  
\hspace{1cm} (A.35)
where $g_t (\xi, k) = \left( \prod_{i=1}^{d_2} L_i \right)^{-1/2} e^{-\pi t^2 |\xi|^2 + |k|^2}$. Furthermore, we note that the dual lattice of $\Delta = \prod_{i=1}^{d_2} L_i^{-1} \mathbb{Z}$ is $\Delta' = \prod_{i=1}^{d_2} L_i \mathbb{Z}$ and the volume of the fundamental domain of $\Delta'$ is $\prod_{i=1}^{d_2} L_i$. Therefore,

$$
\mathcal{F}_\Gamma^\pm \{ g_t \} (x, y) = \mathcal{F}_\mathbb{R}^\pm \left\{ e^{-\pi t^2 |\xi|^2} \right\} (x) \cdot \mathcal{F}_\prod_{i=1}^{d_2} L_i^{-1} \{ e^{-\pi |k|^2} \} (y)
$$

$$
= t^{-d_1} e^{-\pi |x|^2 / t^2} \left( \prod_{i=1}^{d_2} L_i \right)^{-1} \sum_{k \in \prod_{i=1}^{d_2} L_i^{-1} \mathbb{Z}} e^{-\pi |k|^2} e^{2\pi ik \cdot y} = t^{-d_1} e^{-\pi |x|^2 / t^2} \sum_{k \in \prod_{i=1}^{d_2} L_i \mathbb{Z}} e^{-\pi |y+k|^2 / t^2} = t^{-(d_1 + d_2)} e^{-\pi |x|^2 / t^2} \sum_{k \in \prod_{i=1}^{d_2} L_i \mathbb{Z}} e^{-\pi |y+k|^2 / t^2},
$$

(A.36)

where we used the Poisson summation formula for $h(y) = e^{-\pi t^2 |y|^2}$ on lattices to justify the second to last equality. By the fifth item, (A.35) and (A.36), we have

$$
\int_{\Gamma} e^{-\pi t^2 (|\xi|^2 + |k|^2)} e^{\pm 2\pi i \xi \cdot x} e^{\pm 2\pi i k \cdot y} \mathcal{F}_\Gamma^\pm \{ f \} (\xi, k) \, d\xi dk = \int_{\Gamma} f(r, s) \mathcal{F}_\Gamma^\pm \{ g_t \} (r - x, s - y) \, dr ds
$$

$$
= \int_{\Gamma} f(r, s) \mathcal{F}_\Gamma^\pm \{ g_t \} (x - r, y - s) \, dr ds = f \ast \mathcal{F}_\Gamma^\pm \{ g_t \} (x, y).
$$

(A.37)

We note that by the dominated convergence theorem, we have

$$
\lim_{t \to 0^+} \int_{\Gamma} e^{-\pi t^2 (|\xi|^2 + |k|^2)} e^{\pm 2\pi i \xi \cdot x} e^{\pm 2\pi i k \cdot y} \mathcal{F}_\Gamma^\pm \{ f \} (\xi, k) \, d\xi dk = \mathcal{F}_\Gamma^\pm \circ \mathcal{F}_\Gamma^\pm \{ f \} (x, y)
$$

(A.38)

for all $(x, y) \in \Gamma$. On the other hand, we note that $\{ \mathcal{F}_\Gamma^\pm \{ g_t \} \}_{t > 0}$ form an approximate identity on $\Gamma$ since for all $t > 0$, $\mathcal{F}_\Gamma^\pm \{ g_t \} (r, s) \geq 0$ on $\Gamma$ and by (A.36),

$$
\int_{\Gamma} \mathcal{F}_\Gamma^\pm \{ g_t \} (r, s) \, dr ds = t^{-(d_1 + d_2)} \int_{\Gamma} e^{-\pi |r|^2 / t^2} \sum_{k \in \prod_{i=1}^{d_2} L_i \mathbb{Z}} e^{-\pi |s+k|^2 / t^2} \, dr ds
$$

$$
= t^{-(d_1 + d_2)} \int_{\mathbb{R}^{d_1}} e^{-\pi |r|^2 / t^2} \int_{\mathbb{R}^{d_2}} e^{-\pi |s|^2 / t^2} \, ds = 1,
$$

(A.39)

and by (A.39) we also have $\lim_{t \to 0^+} \int_{\Gamma \setminus B(0, \delta)} \mathcal{F}_\Gamma^\pm \{ g_t \} (r, s) \, dr ds \to 0$ for all $\delta > 0$. Therefore, for all $(x, y) \in \Gamma$ we have

$$
\mathcal{F}_\Gamma^\pm \circ \mathcal{F}_\Gamma^\pm \{ f \}(x, y) = \lim_{t \to 0^+} f \ast \mathcal{F}_\Gamma^\pm \{ g_t \} (x, y) = f(x, y).
$$

(A.40)

To prove the analogous inversion formula on $\mathcal{F}(\hat{\Gamma}; \mathbb{C})$, for any given $t > 0$ and $(\xi, k) \in \hat{\Gamma}$ we consider the function

$$
\varphi^\pm (x, y) = \left( \prod_{i=1}^{d_2} L_i \right)^{1/2} e^{-\pi t^2 |x|^2} e^{\pm 2\pi i \xi \cdot x} e^{\pm 2\pi i k \cdot y}.
$$

(A.41)

We note that

$$
\mathcal{F}_\Gamma^\pm \{ \varphi^\pm \} (\eta, m) = \mathcal{F}_\Gamma^\pm \{ h_t \} (\eta - \xi, m - k)
$$

(A.42)

where $h_t (x, y) = \left( \prod_{i=1}^{d_2} L_i \right)^{1/2} e^{-\pi t^2 |x|^2}$, and by a change of variables we have

$$
\mathcal{F}_\Gamma^\pm \{ h_t \} (\xi, k) = \left( \prod_{i=1}^{d_2} L_i \right)^{1/2} \mathcal{F}_{\mathbb{R}^{d_1}} \{ e^{-\pi t^2 |\xi|^2} \} (\xi) \cdot \mathcal{F}_{\prod_{i=1}^{d_2} L_i \mathbb{Z}} \{ \delta_0 \} (k) = t^{-d_1} e^{-\pi |\xi|^2 / t^2} \delta_{0, k},
$$

(A.43)

where $\delta_{0, k}$ is the Kronecker delta. Thus, by the fifth item, we have

$$
\int_{\Gamma} \left( \prod_{i=1}^{d_2} L_i \right)^{1/2} e^{-\pi t^2 |x|^2} e^{\pm 2\pi i \xi \cdot x} e^{\pm 2\pi i k \cdot y} \mathcal{F}_\Gamma^\pm \{ f \} (x, y) \, dx dy = \int_{\Gamma} f(\eta, m) \mathcal{F}_\Gamma^\pm \{ h_t \} (\eta - \xi, m - k) \, d\eta dm
$$

$$
= \int_{\Gamma} f(\eta, m) \mathcal{F}_\Gamma^\pm \{ g_t \} (\xi - \eta, k - m) \, d\eta dm = f \ast \mathcal{F}_\Gamma^\pm \{ h_t \} (\xi, k).
$$

(A.44)
We note that by \((A.42), \{\mathcal{F}_{\Gamma}^{\pm}(h_t)\}_{t>0}\) form an approximate identity on \(\mathbf{\hat{\Gamma}}\), therefore by repeating the same argument on \(\mathcal{S}'(\Gamma; \mathbb{C})\) we find that \(f(\xi, k) = \mathcal{F}_{\Gamma}^{\pm}(\{f\})(\xi, k)\) for all \((\xi, k) \in \mathbf{\hat{\Gamma}}\). This proves the sixth item.

The last item follows immediately from the fourth item and the seventh item.

We then extend the definition of the Fourier transform to the class of tempered distributions.

**Definition A.10.** Suppose \(T_1 \in \mathcal{S}'(\Gamma; \mathbb{C})\) and \(T_2 \in \mathcal{S}'(\mathbf{\hat{\Gamma}}; \mathbb{C})\), where \(\Sigma\) is defined via \((1.43)\).

1. We define the Fourier and inverse Fourier transform of \(T_1\) to be \(\mathcal{F}_{\Gamma}^{\pm}\{T_1\} \in \mathcal{S}'(\mathbf{\hat{\Gamma}}; \mathbb{C})\) by

\[
\langle \mathcal{F}_{\Gamma}^{\pm}\{T_1\}, \varphi \rangle = \langle T_1, \mathcal{F}_{\Gamma}^{\pm}\{\varphi\} \rangle \quad \text{for all } \varphi \in \mathcal{S}(\Gamma; \mathbb{C})
\]

2. We define the Fourier and inverse Fourier transform of \(T_2\) to be \(\mathcal{F}_{\Gamma}^{\pm}\{T_2\} \in \mathcal{S}'(\Gamma; \mathbb{C})\) by

\[
\langle \mathcal{F}_{\Gamma}^{\pm}\{T_2\}, \varphi \rangle = \langle T_2, \mathcal{F}_{\Gamma}^{\pm}\{\varphi\} \rangle \quad \text{for all } \varphi \in \mathcal{S}(\Gamma; \mathbb{C})
\]

We note that by Proposition A.9, the maps \(\mathcal{F}_{\Gamma}^{\pm}: \mathcal{S}'(\Gamma; \mathbb{C}) \to \mathcal{S}'(\mathbf{\hat{\Gamma}}; \mathbb{C})\) and \(\mathcal{F}_{\Gamma}^{\pm}: \mathcal{S}'(\mathbf{\hat{\Gamma}}; \mathbb{C}) \to \mathcal{S}'(\Gamma; \mathbb{C})\) are also isomorphisms.

Throughout the paper we will follow standard notation by denoting the Fourier and inverse Fourier transforms for a Schwartz function \(f \in \mathcal{S}(\Gamma; \mathbb{C})\) or a tempered distribution \(f \in \mathcal{S}'(\Gamma; \mathbb{C})\) by \(\hat{f} = \mathcal{F}_{\Gamma}\{f\}, \check{f} = \mathcal{F}_{\Gamma}^{-1}\{f\}\). Sometimes we will also write \(\mathcal{F}_{\Gamma}[f] = \mathcal{F}_{\Gamma}^{\pm}\{f\}, \mathcal{F}_{\Gamma}^{-1}[f] = \mathcal{F}_{\Gamma}^{\pm}\{-f\}\).

**A.3. Spaces defined via Fourier multipliers.** We note that since \(\mathbf{\hat{\Gamma}}\) and \(\mathbf{\hat{\Gamma}}\) are always subgroups of \(\mathbb{R}^d\), we can extend the reordering map \(P_{g,\Gamma}\) trivially to \(\mathbb{R}^d\). In the next lemma, we show that the restriction of the map \(\sigma_{\Gamma, r}\) on a space \(\mathcal{L}_{\omega}(\Gamma; F)\), where \(\sigma_{\Gamma, r}\) is defined in the third item of Definition A.3, is an isometric isomorphism between \(\mathcal{L}_{\omega}(\Gamma; F)\) and \(\mathcal{L}_{\omega}(g\Gamma; F)\), where \(\mathcal{L}_{\omega}(\Gamma; F)\) is a function space defined in terms of a Fourier multiplier \(\omega : \mathbb{R}^d \to [0, \infty)\) that is invariant under \(P_{g,\Gamma}\).

**Lemma A.11.** Let \(\Gamma\) be defined as in \((1.43)\). Let \(g \in S_{d, g}\), \(g\Gamma\) be defined as in the first item of Definition A.3, and \(P_{g,\Gamma}: P_{g,\Gamma}\) be the reordering maps defined as in \((A.6)\). Let \(\omega : \mathbb{R}^d \to [0, \infty)\) be a Fourier multiplier that is \(L^2\)-almost everywhere such that \(\omega(\cdot) = \omega(P_{g,\Gamma}(\cdot))\). We define the function space \(\mathcal{L}_{\omega}(\Gamma; F)\) via

\[
\mathcal{L}_{\omega}(\Gamma; F) = \left\{ \left. f \in \mathcal{S}'(\Gamma; \mathbb{C}) \right| \hat{f} \in L_{\text{loc}}(\mathbf{\hat{\Gamma}}; \mathbb{C}), \ f = \check{f}, \ \| \sqrt{\omega} f \|_{L^2} < \infty, \right\}, \quad F = \mathbb{R}
\]

\[
\mathcal{L}_{\omega}(\Gamma; F) = \left\{ \left. f \in \mathcal{S}'(\Gamma; \mathbb{C}) \right| \hat{f} \in L_{\text{loc}}(\mathbf{\hat{\Gamma}}; \mathbb{C}), \ \| \sqrt{\omega} f \|_{L^2} < \infty, \right\}, \quad F = \mathbb{C},
\]

Then the map \(\sigma_{\Gamma, r} : \mathcal{L}_{\omega}(\Gamma; F) \to \mathcal{L}_{\omega}(g\Gamma; F)\) defined via \(\sigma_{\Gamma, r} f = f \circ P_{g,\Gamma}^{-1}\) is an isometric isomorphism, with its inverse given by the map \(\sigma_{\Gamma, r}^{-1} : \mathcal{L}_{\omega}(g\Gamma; F) \to \mathcal{L}_{\omega}(\Gamma; F)\) defined via \(\sigma_{\Gamma, r}^{-1} f = f \circ P_{g,\Gamma}\).

**Proof.** We first note that by Lemma A.8, the map \(\sigma_{\Gamma, r} : \mathcal{S}'(\Gamma; \mathbb{C}) \to \mathcal{S}'(g\Gamma; \mathbb{C})\), with its inverse given by \(\sigma_{\Gamma, r}^{-1} : \mathcal{S}'(g\Gamma; \mathbb{C}) \to \mathcal{S}'(\Gamma; \mathbb{C})\). Then for any \(f \in \mathcal{L}_{\omega}(\Gamma; F)\), we have \(\sigma_{\Gamma, r} f \in \mathcal{S}'(g\Gamma; \mathbb{C})\), \(\sigma_{\Gamma, r}^{-1}[\sigma_{\Gamma, r} f] \in L_{\text{loc}}(\mathbf{\hat{\Gamma}}; \mathbb{C})\), and \(\sigma_{\Gamma, r} f = \sigma_{\Gamma, r}^{-1} f\). Furthermore, by \((A.31)\) we have \(\mathcal{F}_{\Gamma}^{\pm}[\sigma_{\Gamma, r} f] = \sigma_{\Gamma, r}^{-1}[\mathcal{F}_{\Gamma}^{\pm}[f]]\). Then

\[
\|\sigma_{\Gamma, r} f\|^2 \mathcal{L}_{\omega}(g\Gamma; F) = \int_{g\Gamma} \omega(\xi)|\mathcal{F}_{\Gamma}^{\pm}[f](P_{g,\Gamma}^{-1}(\xi))|^2 d\xi = \int_\Gamma \omega(\sigma_{\Gamma, r} f)|\mathcal{F}_{\Gamma}^{\pm}[f] f(\xi)|^2 d\xi = \int_\Gamma \omega(\xi)|\mathcal{F}_{\Gamma}^{\pm}[f] f(\xi)|^2 d\xi = \|f\|^2 \mathcal{L}_{\omega}(\Gamma; F)\). \tag{A.48}
\]

This shows that \(\sigma_{\Gamma, r} f \in \mathcal{L}_{\omega}(g\Gamma; F)\) with \(\|\sigma_{\Gamma, r} f\|_{\mathcal{L}_{\omega}(g\Gamma; F)} = \|f\|_{\mathcal{L}_{\omega}(\Gamma; F)}\). Following the exact same calculations as above we also find that for any \(f \in \mathcal{L}_{\omega}(g\Gamma; F)\), we have \(\sigma_{\Gamma, r}^{-1} f = f \circ P_{g,\Gamma} \in \mathcal{L}_{\omega}(\Gamma; F)\) with \(\|\sigma_{\Gamma, r}^{-1} f\|_{\mathcal{L}_{\omega}(\Gamma; F)} = \|f\|_{\mathcal{L}_{\omega}(g\Gamma; F)}\). This shows that \(\sigma_{\Gamma, r} : \mathcal{L}_{\omega}(\Gamma; F) \to \mathcal{L}_{\omega}(g\Gamma; F)\) is an isometric isomorphism with \(\sigma_{\Gamma, r}^{-1}\) as its inverse. □

**Acknowledgement.** The authors would like to thank Benoît Pausader for pointing out the proof for Theorem 2.17.

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