# NONUNIFORM DEPENDENCE ON INITIAL DATA FOR COMPRESSIBLE GAS DYNAMICS: THE CAUCHY PROBLEM ON $\mathbb{R}^{2}$ 

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#### Abstract

The Cauchy problem for the two dimensional compressible Euler equations with data in the Sobolev space $H^{s}\left(\mathbb{R}^{2}\right)$ is known to have a unique solution of the same Sobolev class for a short time, and the data-to-solution map is continuous. We prove that the data-to-solution map on the plane is not uniformly continuous on any bounded subset of Sobolev class functions.


## 1. Introduction

In this paper, we consider the Cauchy problem for the two-dimensional compressible Euler equations with data in the Sobolev space $H^{s}\left(\mathbb{R}^{2}\right)$. The problem can be written in the form

$$
\begin{align*}
& \left\{\begin{array}{l}
\rho_{t}+\rho_{0} u_{x}+(\rho u)_{x}+\rho_{0} v_{y}+(\rho v)_{y}=0 \\
u_{t}+u u_{x}+v u_{y}+h_{x}+\frac{h_{0}+h}{\rho_{0}+\rho} \rho_{x}=0 \\
v_{t}+u v_{x}+v v_{y}+h_{y}+\frac{h_{0}+h}{\rho_{0}+\rho} \rho_{y}=0 \\
h_{t}+u h_{x}+v h_{y}+(\gamma-1)\left(h_{0}+h\right)\left(u_{x}+v_{y}\right)=0
\end{array}\right.  \tag{1.1}\\
& \left.\rho\right|_{t=0}=\phi_{1},\left.\quad u\right|_{t=0}=\phi_{2},\left.\quad v\right|_{t=0}=\phi_{3},\left.\quad h\right|_{t=0}=\phi_{4}, \tag{1.2}
\end{align*}
$$

where $\gamma>1, \rho_{0}>0$ and $h_{0}>0$ are constant. In order to arrive at this from the standard form of the equations for ideal compressible gas dynamics (see for example Majda, [34, pp 3-4]), we have written the density as $\rho_{0}+\rho$ and have replaced the pressure $p$ by a multiple of the internal energy, $h_{0}+h=p /\left(\rho_{0}+\rho\right)$. The velocity components are $u$ and $v$. We have also written the system in nonconservative form, as we are considering only classical solutions in this paper. The purpose of the constants $\rho_{0}$ and $h_{0}$ is to allow us to work with a state variable $U=(\rho, u, v, h)$ whose components lie in the Sobolev space, $H^{s}\left(\mathbb{R}^{2}\right)=H^{s}$, defined as

$$
H^{s}=\left\{f \in \mathcal{S}^{\prime}\left(\mathbb{R}^{2}\right):\left\|\mathcal{F}^{-1}\left(\left(1+|\xi|^{2}\right)^{s / 2} \widehat{f}\right)\right\|_{L^{2}\left(\mathbb{R}^{2}\right)}<\infty\right\} .
$$

Pointwise restrictions on the initial data (see discussion following the statement of Theorem (4) allow us to stay a positive distance from a vacuum state.

Local in time well-posedness in the sense of Hadamard for the system (1.1) (in $d$ space dimensions) is well known when $s>1+d / 2$. The idea of the proof goes back to Gårding [17], Leray, [31], Lax [30] and Kato [25]; a modern version can be found in Taylor's monograph, 40]. For a more detailed exposition of the background and for alternative proofs, see Majda [34] or

[^0]Serre [38]. In particular, if the initial data is in the Sobolev space $H^{s}$, for any $s>1+d / 2$, then there exists a unique solution for some time interval which depends upon the $H^{s}$ norm of the initial data, and the solution depends continuously on the initial condition. In addition, the solution size (in $H^{s}$ ) is bounded by twice the size of the initial condition for some period of time. Classical solutions to the compressible Euler equations do not exist globally in time. Indeed, it has been shown that even for almost constant initial data, there is generally a critical time, $T^{C}$, at which the classical $\left(H^{s}\right)$ solution breaks down [34]. This breakdown is characterized by the formation of shock waves; that is, as $t \nearrow T^{C}$,

$$
\limsup _{t \rightarrow T^{C}}\left\|u_{t}\right\|_{L^{\infty}}+\|\nabla u\|_{L^{\infty}}=\infty
$$

Weak solutions for quasilinear systems in conservation form (the standard form for (1.1)) have been extensively studied in a single space dimension, where there is a complete wellposedness theory for data of small total variation (and in some cases small oscillation). Excellent monographs by the originators of this theory can be found in Bressan [4] and Dafermos [13].

An outstanding open problem in multidimensional hyperbolic conservation laws is to develop a theory of weak solutions for times after the formation of a shock wave. This is an active area of current research. Čanić, Keyfitz, Jegdić and co-authors (for example, [5, 6, 23] and the recent [24]) have looked at self-similar solutions of two-dimensional problems, as have Chen, Feldman and co-authors, [7, 8] for example). There is also interesting work by Shu-Xing Chen, [9] and other papers, and by Elling, see [15] and references there. An intriguing line of research concerns ill-posedness of multi-dimensional problems of the type of (1.1) in spaces other than $H^{s}$; Rauch [37, following Brenner [3], identified key points of this issue, first identified by Littman 32]; and Dafermos [12] and Lopes [33] have followed it up. Yet another question, that concerns the proper definition of weak solutions, is raised by the "wild" weak solutions of De Lellis and Székelyhidi [14. While this does not seem to bear on the question we tackle here, which concerns classical solutions, it is worth mentioning, both as a note about well-posedness, and as evidence of the relationship between the compressible and incompressible gas dynamics equations, which we exploit in this paper.

The compressible Euler equations can be reduced, in the zero Mach number limit, to the incompressible Euler equations (see [34] or [35] for details on the asymptotic analysis),

$$
\begin{align*}
u_{t}+u u_{x}+v u_{y}+p_{x} & =0 \\
v_{t}+u v_{x}+v v_{y}+p_{y} & =0  \tag{1.3}\\
u_{x}+v_{y} & =0
\end{align*}
$$

where $p$ is pressure. Global in time well-posedness is also an important question for the incompressible Euler equations. For a summary of the open questions, we refer the reader to Constantin [10] and Fefferman [16]. For local well-posedness and related results see MajdaBertozzi 35]. Because the incompressible system is not hyperbolic, analysis of the two problems - (1.1) and (1.3) - has proceeded along rather different lines. This paper finds a rather striking connection.

A point of departure for our analysis is the proof of the non-uniform continuity of the data-to-solution map for the incompressible Euler equations recently established by Himonas and Misiołek [21]. In particular, in dimensions two and three they found solutions for periodic data and for Sobolev space data, for which the data-to-solution map was not uniformly continuous. In the non-periodic (full plane) case, their method used a technique of high-low frequency approximate solutions developed by Koch and Tzvetkov [29] for the one-dimensional Benjamin-Ono equation. Our main result is to show that, in a similar way, dependence on the initial data is not better than continuous for classical solutions of the compressible Euler equations. We state our result as follows. (Here we assume the standard restriction on $s, s>d / 2+1$.)

Theorem 1. For $s>2$, the data-to-solution map for the system (1.1) is not uniformly continuous from any bounded subset of $\left(H^{s}\left(\mathbb{R}^{2}\right)\right)^{4}$ to the solution space $C\left([-T, T] ;\left(H^{s}\left(\mathbb{R}^{2}\right)\right)^{4}\right)$.

Our proof of non-uniform dependence of the data-to-solution map uses a method similar to that of [21] and [29]: construction of high-low frequency approximate solutions. We formulate a different way of defining the low frequency terms. In particular, Koch-Tzvetkov and HimonasMisiołek use an $L^{2}$ energy estimate, while we use an energy estimate in $H^{\sigma}, \sigma<s-1$. We are able to do this by sidestepping the construction of some low frequency exact solutions to the compressible Euler equations. The strategy in this paper is to find estimates in the $H^{\sigma}$ norm for $\sigma$ near $s$. We find that the low frequency residual terms actually help to give the desired estimates by allowing for a crucial cancellation. These convenient cancellations, obtained in our construction, simplify technical difficulties created by the more complicated system of equations. The construction of approximate solutions and demonstration of non-uniformity were first carried out in the ideal compressible gas dynamics system (1.1) for periodic data. This result is in the companion paper of Keyfitz and Tığlay [28] along with the description of the flow for the approximate solutions that we use.

Continuity properties of the data-to-solution map for a variety of equations have been studied by many other authors. In particular, the first result of this type was shown by Kato [25] for Burgers' equation, $u_{t}+\left(u^{2}\right)_{x}=0$. Kato showed that the data-to-solution map is not Hölder continuous from any bounded subset of $H^{s}$ to $H^{s}$, when $s>3 / 2$.

The idea of using high-frequency approximate solutions has also been employed extensively in the context of dispersive equations. For example, both Christ, Colliander and Tao [11] and Kenig, Ponce and Vega [27] used a similar method of high-frequency approximate solutions to show ill-posedness of some defocusing dispersive equations. This methodology was also adapted by Himonas and Kenig [19] for the Camassa-Holm (CH) equation on the circle, and by Himonas, Kenig and Misiołek [20] for the CH equation in the non-periodic case. For additional related results concerning the continuity of data-to-solution maps, we refer the reader to Bona and Tzvetkov [2], Holmes [22], Molinet, Saut and Tzvetkov [36] and the references contained therein.

In the next section, we give some preliminary results and notation which we shall use throughout our proof. Section 3 gives the proof of non-uniform dependence.

## 2. Preliminary Results and Notation

This section summarizes background needed in the rest of the paper. The operator $\Lambda^{s} f$ is defined by the formula

$$
\widehat{\Lambda^{s} f}(\xi, \eta)=\left(1+\xi^{2}+\eta^{2}\right)^{s / 2} \widehat{f}(\xi, \eta)
$$

where $f$ is a test function. Here $s$ may be any positive real number; in order to use the standard existence theorems for classical solutions of (1.1), we take $s>d / 2+1=2$. The notation stands for the usual Fourier transform. The Sobolev space $H^{s}$ is a Hilbert space equipped with inner product and norm given by

$$
\|f\|_{s}^{2}=\left\langle\Lambda^{s} f, \Lambda^{s} f\right\rangle_{L^{2}}
$$

We will frequently employ the following Sobolev embedding theorem (see, for instance, Taylor [40, p 272]).
Theorem 2 (Sobolev embedding). If $s>k+1$, then $H^{s}$ is continuously embedded into $C^{k}$ : $H^{s} \hookrightarrow C^{k}$. Specifically,

$$
\begin{equation*}
H^{s} \subset\left\{f \in C^{k}: D^{\alpha} f\left(x_{1}, x_{2}\right) \rightarrow 0 \text { as }\left|\left(x_{1}, x_{2}\right)\right| \rightarrow \infty,|\alpha| \leq k\right\} \tag{2.1}
\end{equation*}
$$

and the inclusion is continuous; for some constant $C(s, k)$ we have

$$
\begin{equation*}
\|f\|_{C^{k}} \doteq \sum_{|\alpha| \leq k}\left\|D^{\alpha} f\right\|_{L^{\infty}} \leq C(s, k)\|f\|_{s} . \tag{2.2}
\end{equation*}
$$

We will also liberally employ the following classical product estimate (see for instance Taylor [40, p 66]).

Lemma 1. If $s>0$ and $f, g \in L^{\infty} \cap H^{s}$, we have the estimate

$$
\|f g\|_{s} \leq C(s)\left[\|f\|_{L^{\infty}}\|g\|_{s}+\|f\|_{s}\|g\|_{L^{\infty}}\right]
$$

This, combined with the Sobolev embedding theorem, implies that $H^{s}$ is a Banach algebra whenever $s>1$; in other words, for $f, g \in H^{s}$, the product $f g \in H^{s}$. Moreover, we have the algebra estimate

$$
\begin{equation*}
\|f g\|_{s} \leq C(s)\|f\|_{s}\|g\|_{s} \tag{2.3}
\end{equation*}
$$

For any test function $f$, the commutator operator $\left[\Lambda^{s}, f\right]$ applied to a test function $g$ is

$$
\begin{equation*}
\left[\Lambda^{s}, f\right] g=\Lambda^{s}(f g)-f \Lambda^{s} g \tag{2.4}
\end{equation*}
$$

The following commutator estimate can be found in Kato and Ponce [26].
Theorem 3 (Kato-Ponce commutator estimate). If $s \geq 0$, and $f \in \operatorname{Lip} \cap H^{s}$ and $g \in L^{\infty} \cap H^{s-1}$, then

$$
\begin{equation*}
\left\|\left[\Lambda^{s}, f\right] g\right\|_{L^{2}} \leq C(s)\left(\|\partial f\|_{L^{\infty}}\left\|\Lambda^{s-1} g\right\|_{L^{2}}+\left\|\Lambda^{s} f\right\|_{L^{2}}\|g\|_{L^{\infty}}\right) . \tag{2.5}
\end{equation*}
$$

In the proof of Lemma 7 we need a simple interpolation estimate:

Proposition 1. If $u \in H^{\tau}$ and $\sigma<s<\tau$ then

$$
\|u\|_{s} \leq\|u\|_{\sigma}^{\alpha}\|u\|_{\tau}^{\beta}, \quad \text { where } \quad \alpha=\frac{\tau-s}{\tau-\sigma}, \quad \beta=\frac{s-\sigma}{\tau-\sigma} .
$$

Proof. We write

$$
\|u\|_{s}^{2}=\int\left(1+|\xi|^{2}\right)^{s}|\hat{u}(\xi)|^{2} d \xi=\int\left[\left(1+|\xi|^{2}\right)^{\sigma}|\hat{u}(\xi)|^{2}\right]^{\alpha}\left[\left(1+|\xi|^{2}\right)^{\tau}|\hat{u}(\xi)|^{2}\right]^{\beta} d \xi
$$

and apply Hölder's inequality with $p=1 / \alpha$ and $q=1 / \beta$.
Finally, owing to the nature of the nonlinearities in (1.1), we need the following reciprocal estimate. It was proved by Kato, [25, Lemma 2.13], for functions in "uniformly local" Sobolev spaces (which generalize our construction of coefficients of the form $\rho_{0}+\rho$ ), and for integer values of $s>2$. We provide a sketch of the proof in the delicate case when $1<s<2$; the larger value of $s$ are straightforward.

Lemma 2. If $s>1, h \in H^{s}\left(\mathbb{R}^{2}\right), g \in H^{s}\left(\mathbb{R}^{2}\right) \cap C^{1}\left(\mathbb{R}^{2}\right)$ and $b>0$ is a constant such that $g+b>\frac{1}{2} b$, then

$$
\left\|\frac{h}{g+b}\right\|_{s} \leq C(s, b)\left(1+\|g\|_{C^{1}}^{s}+\|g\|_{s}^{s}\right)\|h\|_{s} .
$$

Proof. In the case $s=1+\gamma$ with $\gamma \in(0,1)$, the integer parts of the norm satisfy this bound as in Kato, [25, Lemma 2.13]. The fractional portion of the norm (see [39, page 155] for instance, for this form of the Sobolev norm) is

$$
\sup _{|\alpha|=1}\left\|D^{\alpha}\left(h(g+b)^{-1}\right)\right\|_{\dot{H}^{\gamma}}^{2} \equiv \sup _{|\alpha|=1} \int_{\mathbb{R}^{4}} \frac{\left|D^{\alpha}\left(h(g+b)^{-1}\right)(\mathbf{x})-D^{\alpha}\left(h(g+b)^{-1}\right)(\mathbf{y})\right|^{2}}{|\mathbf{x}-\mathbf{y}|^{2 \gamma+2}} d \mathbf{x} d \mathbf{y},
$$

where $\mathbf{x}$ and $\mathbf{y}$ are points in $\mathbb{R}^{2}$ and $\dot{H}^{\gamma}$ is the homogeneous Sobolev space. Consider $D^{\alpha}=\partial_{1}$ (the partial derivative with respect to the first component) so that

$$
\begin{equation*}
\partial_{1}\left(h(g+b)^{-1}\right)=(g+b)^{-1} \partial_{1} h-h(g+b)^{-2} \partial_{1} g . \tag{2.6}
\end{equation*}
$$

Estimating the first term on the right hand side of (2.6) is a straightforward calculation after breaking the integral into the following two pieces

$$
\begin{align*}
\left\|(g+b)^{-1} \partial_{1} h\right\|_{H^{\gamma}}^{2} \leq & 2 \int \frac{1}{|g(\mathbf{x})+b|^{2}} \frac{\left|\partial_{1} h(\mathbf{x})-\partial_{1} h(\mathbf{y})\right|^{2}}{|\mathbf{x}-\mathbf{y}|^{2 \gamma+2}} d \mathbf{x} d \mathbf{y} \\
& +2 \int \frac{|g(\mathbf{x})-g(\mathbf{y})|^{2}}{|\mathbf{x}-\mathbf{y}|^{2 \gamma+2}} \frac{1}{|g(\mathbf{x})+b|^{2}} d \mathbf{x} \frac{\left|\partial_{1} h(\mathbf{y})\right|^{2}}{|g(\mathbf{y})+b|^{2}} d \mathbf{y} \tag{2.7}
\end{align*}
$$

The first integral is bounded by an application of Hölder's inequality, while the second term additionally requires the Sobolev embedding theorem and the following calculus estimate:

$$
\sup _{y \in \mathbb{R}^{2}} \int \frac{|g(\mathbf{x})-g(\mathbf{y})|^{2}}{|\mathbf{x}-\mathbf{y}|^{2 \gamma+2}} d \mathbf{x} \leq C(\gamma)\|g\|_{C^{1}}^{2}
$$

This estimate is obtained by splitting the domain of integration into two pieces, $|x-y|<1$ and $|x-y| \geq 1$, and then applying the mean value theorem. Returning to equation (2.6), the second term on the right hand side is bounded by Lemma 1 and the Sobolev embedding theorem:

$$
\begin{aligned}
\left\|h(g+b)^{-2} \partial_{1} g\right\|_{\gamma} & \leq\left\|(g+b)^{-2} \partial_{1} g\right\|_{\gamma}\|h\|_{L^{\infty}}+\left\|(g+b)^{-2} \partial_{1} g\right\|_{L^{\infty}}\|h\|_{\gamma} \\
& \leq\left(\frac{4}{b^{2}}\|g\|_{s}+\left\|(g+b)^{-2} \partial_{1} g\right\|_{\dot{H}^{\gamma}}\right)\|h\|_{s}+\frac{4}{b^{2}}\|g\|_{C^{1}}\|h\|_{s} .
\end{aligned}
$$

We bound $\left\|(g+b)^{-2} \partial_{1} g\right\|_{\dot{H}^{\gamma}}$ in the same way as $\left\|(g+b)^{-1} \partial_{1} h\right\|_{\dot{H}^{\gamma}}$. The same estimates hold for $\partial_{2}$.

## 3. Proof of Nonuniform Dependence

We write the compressible Euler system (1.1) in the form

$$
\begin{equation*}
U_{t}+A(U) U_{x}+B(U) U_{y}=0, \tag{3.1}
\end{equation*}
$$

with $U=(\rho, u, v, h)^{T}$ and

$$
A(U)=\left(\begin{array}{cccc}
u & \rho_{0}+\rho & 0 & 0 \\
\frac{h_{0}+h}{\rho_{0}+\rho} & u & 0 & 1 \\
0 & 0 & u & 0 \\
0 & (\gamma-1)\left(h_{0}+h\right) & 0 & u
\end{array}\right), \quad B(U)=\left(\begin{array}{cccc}
v & 0 & \rho_{0}+\rho & 0 \\
0 & v & 0 & 0 \\
\frac{h_{0}+h}{\rho_{0}+\rho} & 0 & v & 1 \\
0 & 0 & (\gamma-1)\left(h_{0}+h\right) & v
\end{array}\right) .
$$

3.1. Symmetrized system. The system (3.1) is symmetrizable; that is, it can be written as

$$
A_{0} U_{t}+A_{1}(U) U_{x}+B_{1}(U) U_{y}=0
$$

where the matrices $A_{0}, A_{1}, B_{1}$ are symmetric and $A_{0}$ is positive definite. We can choose

$$
\begin{aligned}
& A_{0}(U)=\left(\begin{array}{cccc}
\frac{h_{0}+h}{\rho_{0}+\rho} & 0 & 0 & 0 \\
0 & \rho_{0}+\rho & 0 & 0 \\
0 & 0 & \rho_{0}+\rho & 0 \\
0 & 0 & 0 & \frac{\rho_{0}+\rho}{(\gamma-1)\left(h_{0}+h\right)}
\end{array}\right) ; \\
& A_{1}(U)=\left(\begin{array}{cccc}
\frac{h_{0}+h}{\rho_{0}+\rho} u & h_{0}+h & 0 & 0 \\
h_{0}+h & \left(\rho_{0}+\rho\right) u & 0 & \rho_{0}+\rho \\
0 & 0 & \left(\rho_{0}+\rho\right) u & 0 \\
0 & \rho_{0}+\rho & 0 & \frac{\left(\rho_{0}+\rho\right) u}{(\gamma-1)\left(h_{0}+h\right)}
\end{array}\right) ; \\
& B_{1}(U)=\left(\begin{array}{cccc}
\frac{h_{0}+h}{\rho_{0}+\rho} v & 0 & h_{0}+h & 0 \\
0 & \left(\rho_{0}+\rho\right) v & 0 & 0 \\
h_{0}+h & 0 & \left(\rho_{0}+\rho\right) v & \rho_{0}+\rho \\
0 & 0 & \rho_{0}+\rho & \frac{\left(\rho_{0}+\rho\right) v}{(\gamma-1)\left(h_{0}+h\right)}
\end{array}\right) .
\end{aligned}
$$

3.2. Approximate solutions. Our strategy, following the template laid out by Himonas and Misiołek [21], is to use two sequences $(\omega= \pm 1)$ of approximate solutions:

$$
U^{\omega, n}=\left(\begin{array}{c}
\rho^{\omega, n}  \tag{3.2}\\
u^{\omega, n} \\
v^{\omega, n} \\
h^{\omega, n}
\end{array}\right)=\left(\begin{array}{c}
0 \\
u_{1}+u_{2} \\
v_{1}+v_{2} \\
0
\end{array}\right)
$$

The approximate solutions contain low frequency functions $u_{1}, v_{1}$ and high frequency functions $u_{2}$ and $v_{2}$. (Our notation suppresses, for clarity, the dependence of the $u_{i}$ and $v_{i}$ on $n$ and $\omega$.) The high frequency functions are defined for a constant $\delta>0$ as

$$
\begin{equation*}
u_{2}=n^{-\delta-s-1} \partial_{y} S, \text { and } v_{2}=-n^{-\delta-s-1} \partial_{x} S \tag{3.3}
\end{equation*}
$$

where $S$ is a stream function, given by

$$
S(x, y, t)=\psi\left(n^{-\delta} x\right) \psi\left(n^{-\delta} y\right) \sin (n y+\omega t)
$$

for a compactly supported nonnegative cutoff function $\psi$ which equals one on $[-2,2]$. Expanding $u_{2}$ and $v_{2}$ gives

$$
\begin{align*}
& u_{2}=n^{-2 \delta-s-1} \psi\left(n^{-\delta} x\right) \psi^{\prime}\left(n^{-\delta} y\right) \sin (n y+\omega t)+n^{-\delta-s} \psi\left(n^{-\delta} x\right) \psi\left(n^{-\delta} y\right) \cos (n y+\omega t) \\
& v_{2}=-n^{-2 \delta-s-1} \psi^{\prime}\left(n^{-\delta} x\right) \psi\left(n^{-\delta} y\right) \sin (n y+\omega t) \tag{3.4}
\end{align*}
$$

The low frequency functions, $u_{1}$ and $v_{1}$, are

$$
\begin{equation*}
u_{1}=\omega n^{-1} \varphi_{1}\left(n^{-\delta} x\right) \varphi_{2}^{\prime}\left(n^{-\delta} y\right), \text { and } v_{1}=-\omega n^{-1} \varphi_{1}^{\prime}\left(n^{-\delta} x\right) \varphi_{2}\left(n^{-\delta} y\right) \tag{3.5}
\end{equation*}
$$

where $\varphi_{1}^{\prime}$ and $\varphi_{2}$ are also smooth compactly supported functions; $\varphi_{1}^{\prime}$ is identically 1 on the support of $\psi^{\prime}$ and $\varphi_{2} \equiv 1$ on $\operatorname{supp} \psi$. The following cancellation holds.

Lemma 3. For $u$ and $v$ defined in (3.2), (3.4), and (3.5), we have $\partial_{x} u^{\omega, n}+\partial_{y} v^{\omega, n}=0$.
Proof. We have

$$
u_{1, x}=\frac{\omega}{n^{1+\delta}} \phi_{1}^{\prime}\left(\frac{x}{n^{\delta}}\right) \phi_{2}^{\prime}\left(\frac{y}{n^{\delta}}\right)=-v_{1, y}
$$

Considering the high frequency terms, we see from (3.3)

$$
\partial_{x} u_{2}+\partial_{y} v_{2}=n^{-\delta-s-1} \partial_{x} \partial_{y} S-n^{-\delta-s-1} \partial_{y} \partial_{x} S=0
$$

As a result of the definition, the approximate solutions satisfy

$$
U_{t}^{\omega, n}+A\left(U^{\omega, n}\right) U_{x}^{\omega, n}+B\left(U^{\omega, n}\right) U_{y}^{\omega, n}=R
$$

where

$$
R=\left(\begin{array}{c}
0  \tag{3.6}\\
R_{2} \\
R_{3} \\
0
\end{array}\right)=\left(\begin{array}{c}
0 \\
\partial_{t} u^{\omega, n}+u^{\omega, n} \partial_{x} u^{\omega, n}+v^{\omega, n} \partial_{y} u^{\omega, n} \\
\partial_{t} v^{\omega, n}+u^{\omega, n} \partial_{x} v^{\omega, n}+v^{\omega, n} \partial_{y} v^{\omega, n} \\
0
\end{array}\right)
$$

Denote the inner product of two vectors, $V$ and $W$, by $\langle V, W\rangle=\sum\left\langle V_{i}, W_{i}\right\rangle_{L^{2}}$, and for any vector $U$ denote

$$
\begin{equation*}
\|U\|_{\sigma}^{2}=\left\langle\Lambda^{\sigma} U, \Lambda^{\sigma} U\right\rangle=\|\rho\|_{\sigma}^{2}+\|u\|_{\sigma}^{2}+\|v\|_{\sigma}^{2}+\|h\|_{\sigma}^{2} . \tag{3.7}
\end{equation*}
$$

Let $U_{\omega, n}=\left(\rho_{\omega, n}, u_{\omega, n}, v_{\omega, n}, h_{\omega, n}\right)^{T}$ be the actual solution to the Cauchy problem corresponding to equation (3.1), with the same data:

$$
U_{\omega, n}(x, y, 0)=U^{\omega, n}(x, y, 0)=\left(0, u_{1}(x, y, 0)+u_{2}(x, y, 0), v_{1}(x, y, 0)+v_{2}(x, y, 0), 0\right),
$$

again with $\omega=1$ or -1 .
The actual solution is unique, and exists on a time interval which depends only upon the size (in the $H^{s}$ norm) of the initial data and on its distance from the boundary of the region of state space (called $G$ in the statement below) where the system is hyperbolic. We quote the following theorem found in 34.

Theorem 4 (Majda, Theorem 2.1). Assume $U(\cdot, 0)=U_{0} \in H^{s}, s>d / 2+1$ and $U_{0}(x) \in G_{1}$, $\bar{G}_{1} \subset \subset G$. Then there is a time interval $[0, T]$ with $T>0$, so that the equations (1.1) have $a$ unique classical solution $U \in C\left([0, T] ;\left(H^{s}\right)^{4}\right) \cap C^{1}\left([0, T] ;\left(H^{s-1}\right)^{4}\right)$, and $U(x, t) \in G_{2}, \bar{G}_{2} \subset \subset G$, for $(x, t) \in \mathbb{R}^{2} \times[0, T]$; here $T=T\left(\left\|U_{0}\right\|_{s}, G_{1}\right)$.

In our coordinate system, $G=\left\{\rho>-\rho_{0}\right\}$. Having specified values for $\rho_{0}>0$ and $h_{0}$, we might choose, for example, data to lie in a bounded set

$$
G_{1}=\left\{-\frac{1}{4} \rho_{0}<\rho(\cdot, 0)<M_{\rho},|u(\cdot, 0)|<M_{u},|v(\cdot, 0)|<M_{u},-\frac{1}{4} h_{0}<h(\cdot, 0)<M_{h}\right\},
$$

and then take $G_{2}$ to be

$$
G_{2}=\left\{-\frac{1}{2} \rho_{0}<\rho(\cdot, 0)<2 M_{\rho},|u(\cdot, 0)|<2 M_{u},|v(\cdot, 0)|<M_{u},-\frac{1}{2} h_{0}<h(\cdot, 0)<2 M_{h}\right\},
$$

where $M_{\rho}, M_{u}$ and $M_{h}$ are positive numbers. The significant bound, which we need throughout, is the lower bound on $\rho$ in $G_{2}$. Additionally, continuous dependence on the data yields the following $H^{s}$ solution size estimate.

Theorem 5 ([34], Theorem 2.2). Assume $\left.U\right|_{t=0} \in H^{s}, s>2$ and $\left.U\right|_{t=0} \in G_{1}$. There exists a $T^{*}, 0<T^{*} \leq T$ such that

$$
\sup _{t \in\left[0, T^{*}\right]}\|U\|_{s} \leq 2\left\|\left.U\right|_{t=0}\right\|_{s}
$$

In what follows, we take $T^{*}$ to be the value given by this theorem.
We obtain the proof of Theorem by showing the following properties of the corresponding solutions:
(1) Boundedness of initial data (proved in Section 3.3):

$$
\begin{equation*}
\left\|U^{\omega, n}(\cdot, 0)\right\|_{s}=\left\|U_{\omega, n}(\cdot, 0)\right\|_{s}=\left\|\left(\rho_{\omega, n}(0), u_{\omega, n}(0), v_{\omega, n}(0), h_{\omega, n}(0)\right)\right\|_{s} \leq C, \tag{3.8}
\end{equation*}
$$

uniformly in $n$.
(2) Convergence of initial data (Section 3.3): for $\delta<1 / 2$,

$$
\begin{equation*}
\lim _{n \rightarrow \infty}\left\|U_{1, n}(\cdot, 0)-U_{-1, n}(\cdot, 0)\right\|_{s}=0 \tag{3.9}
\end{equation*}
$$

(3) Uniformity of approximation of $U^{\omega, n}$ to actual solution $U_{\omega, n}$ (Section 3.4):

$$
\begin{equation*}
\left\|U_{\omega, n}(\cdot, t)-U^{\omega, n}(\cdot, t)\right\|_{s} \leq C n^{-\varepsilon}, \quad 0<t<T^{*} \tag{3.10}
\end{equation*}
$$

for some $\varepsilon>0$.
(4) Non-uniformity of divergence of $U_{1, n}$ and $U_{-1, n}$ from each other in time (Section 3.7):

$$
\begin{equation*}
\left\|U_{1, n}(\cdot, t)-U_{-1, n}(\cdot, t)\right\|_{s}>|\sin (t)|, \quad 0<t \tag{3.11}
\end{equation*}
$$

The following estimates can be found in the appendix of [21].
Lemma 4. Let $\sigma \geq 0, \delta>0$ and $n \gg 1$. For any Schwarz function $\psi \in \mathcal{S}(\mathbb{R})$ we have

$$
\begin{equation*}
\left\|\psi\left(n^{-\delta} \cdot\right)\right\|_{H^{\sigma}(\mathbb{R})} \leq n^{\delta / 2}\|\psi\|_{H^{\sigma}(\mathbb{R})} \tag{3.12}
\end{equation*}
$$

For any constant $a \in \mathbb{R}$ we have

$$
\begin{equation*}
\left\|\psi\left(n^{-\delta} \cdot\right) \sin (n \cdot+a)\right\|_{H^{\sigma}(\mathbb{R})}+\left\|\psi\left(n^{-\delta} \cdot\right) \cos (n \cdot+a)\right\|_{H^{\sigma}(\mathbb{R})} \approx n^{\sigma+\delta / 2}\|\psi\|_{L^{2}(\mathbb{R})} . \tag{3.13}
\end{equation*}
$$

The notation $\approx$ means that the expression on the left is bounded above and below by constants independent of $\sigma, \delta$ and $n$. Note that the $L^{2}$ bound implies an $H^{\sigma}$ bound. From this lemma we obtain bounds on the approximate solutions:

Lemma 5. For $s-2<\sigma<s-1$ and $0<\delta<1$, we have

$$
\left\|U^{\omega, n}\right\|_{\sigma+1} \leq C n^{\sigma-s+1}
$$

where $C$ depends on the norms of the functions $\psi, \phi_{1}$ and $\phi_{2}$.
Proof. The nonzero terms in $U^{\omega, n}$ are $u_{1}, u_{2}, v_{1}$ and $v_{2}$. Since $u_{1}$ and $v_{1}$ are products (in $x$ and $y$ ) of terms of the form $\psi\left(n^{-\delta} \cdot\right)$, we have, from Lemma 4, and with $C$ a generic constant,

$$
\left\|u_{1}\right\|_{\sigma+1}=n^{-1}\left\|\phi_{1}\left(n^{-\delta} x\right)\right\|_{\sigma+1}\left\|\phi_{2}^{\prime}\left(n^{-\delta} y\right)\right\|_{\sigma+1} \leq n^{-1} n^{\delta / 2}\left\|\phi_{1}\right\|_{\sigma+1} n^{\delta / 2}\left\|\phi_{2}^{\prime}\right\|_{\sigma+1}=C n^{-1+\delta}
$$

A similar bound holds for $v_{1}$, which has the same structure. Note that these bounds are valid for any $\sigma$. On the other hand,

$$
\begin{aligned}
\left\|u_{2}\right\|_{\sigma+1} \leq & n^{-2 \delta-s-1}\left\|\psi\left(n^{-\delta} x\right)\right\|_{\sigma+1}\left\|\psi^{\prime}\left(n^{-\delta} y\right) \sin (n y+\omega t)\right\|_{\sigma+1} \\
& \quad+n^{-\delta-s}\left\|\psi\left(n^{-\delta} x\right)\right\|_{\sigma+1}\left\|\psi\left(n^{-\delta} y\right) \cos (n y+\omega t)\right\|_{\sigma+1} \\
& \leq n^{-2 \delta-s-1} n^{\delta / 2}\|\psi\|_{\sigma+1} n^{\sigma+1+\delta / 2}\left\|\psi^{\prime}\right\|_{\sigma+1}+n^{-\delta-s} n^{\delta / 2}\|\psi\|_{\sigma+1} n^{\sigma+1+\delta / 2}\|\psi\|_{\sigma+1} \\
\leq & C n^{-\delta+\sigma-s}+C n^{\sigma-s+1}
\end{aligned}
$$

while $v_{2}$, which has the structure of the first term in $u_{2}$, satisfies

$$
\left\|v_{2}\right\|_{\sigma+1} \leq C n^{-\delta+\sigma-s} .
$$

Now, $\delta$ is a positive number and $\sigma<s-1$, so all the exponents of $n$ are negative if $\delta<1$. To bound the low frequency terms by the high frequency terms we need $-1+\delta<\sigma-s+1$, or
$\delta<\sigma-s+2$, and provided $\sigma>s-2$, as we have assumed, it is possible to achieve this with $\delta>0$.

Lemma 5 implies a bound on the actual solution, using Theorem 5
Corollary 1. If $t \leq T^{*}$, then

$$
\left\|U_{\omega, n}\right\|_{\sigma+1} \leq C n^{\sigma+1-s}
$$

where $T^{*}$, as in Theorem [5, is the time to doubling of the initial norm, and $\sigma<s-1$.
3.3. (1) Boundedness and (2) Convergence of the Initial Data. From Lemma4, we have $\left\|U^{\omega, n}(\cdot, 0)\right\|_{s} \leq C n^{-1+\delta}+C$. For any $\delta$ with $0<\delta<1$ we have $\left\|U^{\omega, n}(0)\right\|$ bounded, uniformly in $n$.

To see that the difference in the initial data for $\omega= \pm 1$ converges to zero in $H^{s}$, we calculate $U^{1, n}-U^{-1, n}$ at $t=0$, noting that the oscillatory terms cancel at $t=0$, leaving only

$$
\left(\begin{array}{c}
\rho_{1, n}(0)-\rho_{-1, n}(0) \\
u_{1, n}(0)-u_{-1, n}(0) \\
v_{1, n}(0)-v_{-1, n}(0) \\
h_{1, n}(0)-h_{-1, n}(0)
\end{array}\right)=\left(\begin{array}{c}
0 \\
2 n^{-1} \varphi_{1}\left(n^{-\delta} x\right) \varphi_{2}^{\prime}\left(n^{-\delta} y\right) \\
2 n^{-1} \varphi_{1}^{\prime}\left(n^{-\delta} x\right) \varphi_{2}\left(n^{-\delta} y\right) \\
0
\end{array}\right) .
$$

This tends to zero in $H^{s}$ by the first estimate in Lemma 团 for any $\delta \in(0,1)$, as in the first estimate in the proof of Lemma 5 ,
3.4. (3) Uniformity of the Approximation. In this subsection, we denote the actual solutions by $U$. Let $\mathcal{E}=U-U^{\omega, n}=(E, F, G, H)^{T}$ be the error, the difference between the actual and approximate solutions. The main result of this section is

Theorem 6. For $\max \{1, s-2\}<\sigma<s-1, \mathcal{E}$ satisfies

$$
\frac{d}{d t}\|\mathcal{E}\|_{\sigma} \lesssim n^{\sigma+1-s}\|\mathcal{E}\|_{\sigma}+n^{\delta-2}
$$

Furthermore, we have on the time interval of existence

$$
\|\mathcal{E}\|_{\sigma} \lesssim n^{\delta-3+s-\sigma} .
$$

Proof. An equation for the error (the symmetric form of the system is useful here) is

$$
\begin{equation*}
A_{0}\left(U^{\omega, n}\right) \mathcal{E}_{t}+A_{1}\left(U^{\omega, n}\right) \mathcal{E}_{x}+B_{1}\left(U^{\omega, n}\right) \mathcal{E}_{y}+C(U) \mathcal{E}+A_{0}\left(U^{\omega, n}\right) R=0 \tag{3.14}
\end{equation*}
$$

where

$$
\begin{aligned}
C(U) & =A_{0}\left(U^{\omega, n}\right)\left(\begin{array}{cccc}
u_{x}+v_{y} & \rho_{x} & \rho_{y} & 0 \\
-\frac{h_{0} \rho_{x}}{\rho \rho_{x}+\rho_{0}^{2}} & u_{x} & u_{y} & \frac{\rho_{x}}{\rho+\rho_{0}} \\
-\frac{h_{0} \rho_{y}}{\rho \rho_{0}+\rho_{0}^{2}} & v_{x} & v_{y} & \frac{\rho_{y}}{\rho+\rho_{0}} \\
0 & h_{x} & h_{y} & (\gamma-1)\left(u_{x}+v_{y}\right)
\end{array}\right) \\
& =\left(\begin{array}{cccc}
\frac{\left(u_{x}+v_{y}\right) h_{0}}{\rho_{0}} & \frac{h_{0}}{\rho_{0}} \rho_{x} & \frac{h_{0}}{\rho_{0}} \rho_{y} & 0 \\
-\frac{h_{0} \rho_{x}}{\rho+\rho_{0}} & \rho_{0} u_{x} & \rho_{0} u_{y} & \frac{\rho_{0} \rho_{x}}{\rho+\rho_{0}} \\
-\frac{h_{0} \rho_{y}}{\rho+\rho_{0}} & \rho_{0} v_{x} & \rho_{0} v_{y} & \frac{\rho_{0} \rho_{y}}{\rho+\rho_{0}} \\
0 & \frac{\rho_{0}}{(\gamma-1) h_{0}} h_{x} & \frac{\rho_{0}}{(\gamma-1) h_{0}} h_{y} & \frac{\rho_{0}\left(u_{x}+v_{y}\right)}{h_{0}}
\end{array}\right) .
\end{aligned}
$$

We write $C(U) \mathcal{E}$ as

$$
C(U) \mathcal{E}=\left(\begin{array}{c}
\frac{\left(u_{x}+v_{y}\right) h_{0}}{\rho_{0}} E+\frac{h_{0}}{\rho_{0}} \rho_{x} F+\frac{h_{0}}{\rho_{0}} \rho_{y} G \\
-\frac{h_{0} \rho_{x}}{\rho_{0} \rho_{0}} E+\rho_{0} u_{x} F+\rho_{0} u_{y} G+\frac{\rho_{0} \rho_{x}}{\rho+\rho_{0}} H \\
-\frac{h_{0} \rho_{y}}{\rho+\rho_{0}} E+\rho_{0} v_{x} F+\rho_{0} v_{y} G+\frac{\rho_{0} \rho_{y}}{\rho+\rho_{0}} H \\
\frac{\rho_{0}}{(\gamma-1) h_{0}} h_{x} F+\frac{\rho_{0}}{(\gamma-1) h_{0}} h_{y} G+\frac{\rho_{0}\left(u_{x}+v_{y}\right)}{h_{0}} H
\end{array}\right)=\left(\begin{array}{l}
C_{1} \\
C_{2} \\
C_{3} \\
C_{4}
\end{array}\right) .
$$

We apply the operator $\Lambda^{\sigma}$, where $\sigma>1$ and $s-2<\sigma<s-1$, to the left hand side of (3.14) and then take the inner product with $\Lambda^{\sigma} \mathcal{E}$ to obtain

$$
\begin{align*}
\left\langle\Lambda^{\sigma}\left(A_{0}\left(U^{\omega, n}\right) \mathcal{E}_{t}\right), \Lambda^{\sigma} \mathcal{E}\right\rangle= & -\left\langle\Lambda^{\sigma}(C(U) \mathcal{E}), \Lambda^{\sigma} \mathcal{E}\right\rangle  \tag{3.15}\\
& -\left\langle\Lambda^{\sigma}\left(\operatorname{diag}\left(A_{1}\left(U^{\omega, n}\right)\right) \mathcal{E}_{x}+\operatorname{diag}\left(B_{1}\left(U^{\omega, n}\right)\right) \mathcal{E}_{y}\right), \Lambda^{\sigma} \mathcal{E}\right\rangle  \tag{3.16}\\
& -\left\langle\Lambda^{\sigma}\left(A_{R}\left(U^{\omega, n}\right) \mathcal{E}_{x}+B_{R}\left(U^{\omega, n}\right) \mathcal{E}_{y}\right), \Lambda^{\sigma} \mathcal{E}\right\rangle  \tag{3.17}\\
& -\left\langle\Lambda^{\sigma} A_{0}\left(U^{\omega, n}\right) R, \Lambda^{\sigma} \mathcal{E}\right\rangle, \tag{3.18}
\end{align*}
$$

where $\operatorname{diag}(A)$ denotes the diagonal part of a matrix $A$ and $A_{R}=A-\operatorname{diag}(A)$. We now bound the terms on the right hand side.
Estimate for (3.15). We have

$$
\begin{equation*}
\left\langle\Lambda^{\sigma} C(U) \mathcal{E}, \Lambda^{\sigma} \mathcal{E}\right\rangle=\left\langle\Lambda^{\sigma} C_{1}, \Lambda^{\sigma} E\right\rangle+\left\langle\Lambda^{\sigma} C_{2}, \Lambda^{\sigma} F\right\rangle+\left\langle\Lambda^{\sigma} C_{3}, \Lambda^{\sigma} G\right\rangle+\left\langle\Lambda^{\sigma} C_{4}, \Lambda^{\sigma} H\right\rangle . \tag{3.19}
\end{equation*}
$$

These terms are all estimated in a similar way. The Cauchy-Schwarz inequality yields

$$
\left|\left\langle\Lambda^{\sigma} C_{1}, \Lambda^{\sigma} E\right\rangle\right| \leq\left\|\Lambda^{\sigma} C_{1}\right\|_{L_{2}}\left\|\Lambda^{\sigma} E\right\|_{L_{2}}=\left\|C_{1}\right\|_{\sigma}\|E\|_{\sigma} \leq\left\|C_{1}\right\|_{\sigma}\|\mathcal{E}\|_{\sigma}
$$

and so on for the other three terms. To estimate $\left\|C_{i}\right\|_{\sigma}$, we note that all of the $C_{i}$ are of the form

$$
\begin{equation*}
C_{i}=a_{1} E+a_{2} F+a_{3} G+a_{4} H, \tag{3.20}
\end{equation*}
$$

where, up to constant multiples, each $a_{j}$ consists of a derivative, or sum of derivatives, of components of $U$, in some cases divided by $\rho+\rho_{0}$. So, taking $C_{2}$ as an example, and looking at the first summand, we have

$$
\begin{equation*}
\left\|a_{1} E\right\|_{\sigma}=h_{0}\left\|\frac{\rho_{x}}{\rho+\rho_{0}} E\right\|_{\sigma} \leq h_{0}\left\|\frac{\rho_{x}}{\rho+\rho_{0}}\right\|_{\sigma}\|E\|_{\sigma} \leq h_{0}\left\|\frac{\rho_{x}}{\rho+\rho_{0}}\right\|_{\sigma}\|\mathcal{E}\|_{\sigma}, \tag{3.21}
\end{equation*}
$$

where we have used the algebra property, Lemma 2.3. Now we use Lemma 2 to obtain

$$
\left\|\frac{\rho_{x}}{\rho+\rho_{0}}\right\|_{\sigma} \leq C\left(\sigma, \rho_{0}\right)\left(1+\|\rho\|_{\sigma}^{\sigma}\right)\left\|\rho_{x}\right\|_{\sigma} \leq C\|\rho\|_{\sigma+1}
$$

since $\left\|\rho_{x}\right\|_{\sigma} \leq\|\rho\|_{\sigma+1}$, and from Corollary Twe can absorb all the other factors into a constant that depends on $\sigma, \rho_{0}$ and on the $H^{\sigma}$ bound on $\rho$. Finally, estimating $\|\rho\|_{\sigma+1} \leq C n^{\sigma-s+1}$ as in Corollary (1, and treating the other terms in (3.20) in the same way as (3.21), we have

$$
\left\langle\Lambda^{\sigma} C(U) \mathcal{E}, \Lambda^{\sigma} \mathcal{E}\right\rangle \leq C n^{\sigma-s+1}\|\mathcal{E}\|_{\sigma}^{2},
$$

with a constant $C$ that depends upon $\rho_{0}, h_{0}, \gamma$ and $\sigma$. (Since $\|U\|_{\sigma+1}$ decreases with $n$, we can eliminate the dependence of the constant on $U$.)

Estimate of (3.16). We have (up to a sign)

$$
\begin{aligned}
(3.16) & =\left\langle\Lambda^{\sigma}\left(\operatorname{diag}\left(A_{1}\left(U^{\omega, n}\right)\right) \mathcal{E}_{x}+\operatorname{diag}\left(B_{1}\left(U^{\omega, n}\right)\right) \mathcal{E}_{y}\right), \Lambda^{\sigma} \mathcal{E}\right\rangle \\
& =\left\langle\Lambda^{\sigma}\left(\frac{h_{0}}{\rho_{0}} u^{\omega, n} E_{x}+\frac{h_{0}}{\rho_{0}} v^{\omega, n} E_{y}\right), \Lambda^{\sigma} E\right\rangle_{L^{2}}+\left\langle\Lambda^{\sigma}\left(\rho_{0} u^{\omega, n} F_{x}+\rho_{0} v^{\omega, n} F_{y}\right), \Lambda^{\sigma} F\right\rangle_{L^{2}} \\
& +\left\langle\Lambda^{\sigma}\left(\rho_{0} u^{\omega, n} G_{x}+\rho_{0} v^{\omega, n} G_{y}\right), \Lambda^{\sigma} G\right\rangle_{L^{2}}+\left\langle\Lambda^{\sigma}\left(\frac{\rho_{0} u^{\omega, n}}{(\gamma-1) h_{0}} H_{x}+\frac{\rho_{0} v^{\omega, n}}{(\gamma-1) h_{0}} H_{y}\right), \Lambda^{\sigma} H\right\rangle_{L^{2}} .
\end{aligned}
$$

The eight terms in this expression are similar to each other; we show how the first is estimated. Ignoring the constant $h_{0} / \rho_{0}$, consider

$$
I_{1} \equiv\left\langle\Lambda^{\sigma}\left(u^{\omega, n} E_{x}\right), \Lambda^{\sigma} E\right\rangle_{L^{2}}=\int_{\mathbb{R}^{2}} \Lambda^{\sigma}\left(u^{\omega, n} E_{x}\right) \Lambda^{\sigma} E d x d y
$$

This can be written as (recall equation (2.4) for the definition of the commutator)

$$
I_{1}=\int_{\mathbb{R}^{2}}\left(\left[\Lambda^{\sigma},\left(u^{\omega, n}\right)\right] E_{x}+\left(u^{\omega, n}\right) \Lambda^{\sigma} \partial_{x} E\right) \Lambda^{\sigma} E d x d y
$$

We split this integral into two pieces, and apply the Cauchy-Schwarz estimate to the first term, to obtain

$$
I_{1} \leq\left\|\left[\Lambda^{\sigma},\left(u^{\omega, n}\right)\right] E_{x}\right\|_{L^{2}}\|E\|_{\sigma}+\left|\int_{\mathbb{R}^{2}} u^{\omega, n} \Lambda^{\sigma} \partial_{x} E \Lambda^{\sigma} E d x d y\right| .
$$

Now, the Kato-Ponce commutator estimate, (2.5), applied to the first factor gives

$$
\begin{aligned}
\left\|\left[\Lambda^{\sigma},\left(u^{\omega, n}\right)\right] E_{x}\right\|_{L^{2}} & \leq C(\sigma)\left(\left\|u_{x}^{\omega, n}\right\|_{L^{\infty}}\left\|\Lambda^{\sigma-1} E_{x}\right\|_{L_{2}}+\left\|\Lambda^{\sigma} u^{\omega, n}\right\|_{L_{2}}\left\|E_{x}\right\|_{L^{\infty}}\right) \\
& \leq C(\sigma)\left(\left\|u^{\omega, n}\right\|_{\sigma+1}\left\|\Lambda^{\sigma} E\right\|_{L_{2}}+\left\|u^{\omega, n}\right\|_{\sigma}\left\|E_{x}\right\|_{L^{\infty}}\right),
\end{aligned}
$$

using the Sobolev embedding theorem, Theorem 2, which applies here since $\sigma+1>2$. Since we can replace $\left\|u^{\omega, n}\right\|_{\sigma}$ by $\left\|u^{\omega, n}\right\|_{\sigma+1}$, and, using the same Sobolev embedding, replace $\left\|E_{x}\right\|_{L^{\infty}}$ by $\|E\|_{\sigma}$, we obtain

$$
\left\|\left[\Lambda^{\sigma},\left(u^{\omega, n}\right)\right] E_{x}\right\|_{L^{2}} \leq C(s)\left\|u^{\omega, n}\right\|_{\sigma+1}\|E\|_{\sigma}^{2} .
$$

For the second term, integration by parts followed by Hölder's inequality yields

$$
\begin{aligned}
\left|\int_{\mathbb{R}^{2}} u^{\omega, n} \Lambda^{\sigma} \partial_{x} E \Lambda^{\sigma} E d x d y\right| & =\left|\frac{1}{2} \int_{\mathbb{R}^{2}} u_{x}^{\omega, n}\left(\Lambda^{\sigma} E\right)^{2} d x d y\right| \leq \frac{\left\|u_{x}^{\omega, n}\right\|_{L^{\infty}}}{2} \int_{\mathbb{R}^{2}}\left(\Lambda^{\sigma} E\right)^{2} d x d y \\
& =\frac{1}{2}\left\|u_{x}^{\omega, n}\right\|_{L^{\infty}}\|E\|_{\sigma}^{2}
\end{aligned}
$$

and we get a bound similar to the first term, so that

$$
I_{1} \leq C\left\|u^{\omega, n}\right\|_{\sigma+1}\|E\|_{\sigma}^{2} \leq C n^{\sigma+1-s}\|E\|_{\sigma}^{2},
$$

from Corollary 1, with $C=C(\sigma)$. Proceeding the same way with the other seven terms, we obtain

$$
\left|\left\langle\Lambda^{\sigma}\left(\operatorname{diag}\left(A_{1}\left(U^{\omega, n}\right)\right) \mathcal{E}_{x}+\operatorname{diag}\left(B_{1}\left(U^{\omega, n}\right)\right) \mathcal{E}_{y}\right), \Lambda^{\sigma} \mathcal{E}\right\rangle\right| \leq C n^{\sigma+1-s}\|\mathcal{E}\|_{\sigma}^{2},
$$

with the constant depending on on $\rho_{0}, h_{0}, \gamma$ and $\sigma$.
Estimate of (3.17). Inserting the off-diagonal elements of $A_{1}$ and $B_{1}$ from Section 3.1 (note that they are all constant since $h^{\omega, n}=0=\rho^{\omega, n}$ ), we have

$$
\begin{aligned}
-(\sqrt[3.17)]{ } & =\left\langle\Lambda^{\sigma}\left(h_{0} F_{x}+h_{0} G_{y}\right), \Lambda^{\sigma} E\right\rangle_{L^{2}}+\left\langle\Lambda^{\sigma}\left(h_{0} E_{x}+\rho_{0} H_{x}\right), \Lambda^{\sigma} F\right\rangle_{L^{2}} \\
& +\left\langle\Lambda^{\sigma}\left(h_{0} E_{y}+\rho_{0} H_{y}\right), \Lambda^{\sigma} G\right\rangle_{L^{2}}+\left\langle\Lambda^{\sigma}\left(\rho_{0} F_{x}+\rho_{0} G_{y}\right), \Lambda^{\sigma} H\right\rangle_{L^{2}} .
\end{aligned}
$$

Writing the above as an integral and rearranging terms gives

$$
\begin{aligned}
-(\sqrt[3.17)]{ }= & \int_{\mathbb{R}^{2}} h_{0}\left(\partial_{x}\left(\Lambda^{\sigma} E\right) \Lambda^{\sigma} F+\Lambda^{\sigma} E \partial_{x}\left(\Lambda^{\sigma} F\right)+\partial_{y}\left(\Lambda^{\sigma} E\right) \Lambda^{\sigma} G+\Lambda^{\sigma} E \partial_{y}\left(\Lambda^{\sigma} G\right)\right) d x d y \\
& +\int_{\mathbb{R}^{2}} \rho_{0}\left(\partial_{x}\left(\Lambda^{\sigma} F\right) \Lambda^{\sigma} H+\Lambda^{\sigma} F \partial_{x}\left(\Lambda^{\sigma} H\right)+\partial_{y}\left(\Lambda^{\sigma} G\right) \Lambda^{\sigma} H+\Lambda^{\sigma} G \partial_{y}\left(\Lambda^{\sigma} H\right)\right) d x d y \\
= & h_{0} \int_{\mathbb{R}^{2}} \partial_{x}\left(\Lambda^{\sigma} E \Lambda^{\sigma} F\right)+\partial_{y}\left(\Lambda^{\sigma} E \Lambda^{\sigma} G\right) d x d y \\
& +\rho_{0} \int_{\mathbb{R}^{2}} \partial_{x}\left(\Lambda^{\sigma} F \Lambda^{\sigma} H\right)+\partial_{y}\left(\Lambda^{\sigma} G \Lambda^{\sigma} H\right) d x d y
\end{aligned}
$$

and therefore they all integrate to zero.
Estimate of (3.18). Since $A_{0}$ is diagonal and $A_{0}\left(U^{\omega, n}\right)$ is constant, we have

$$
-(3.18)=\left\langle\Lambda^{\sigma} A_{0}\left(U^{\omega, n}\right) R, \Lambda^{\sigma} \mathcal{E}\right\rangle=\left\langle\rho_{0} \Lambda^{\sigma} R_{2}, \Lambda^{\sigma} F\right\rangle+\left\langle\rho_{0} \Lambda^{\sigma} R_{3}, \Lambda^{\sigma} G\right\rangle
$$

The Cauchy-Schwarz inequality yields

$$
|(\sqrt{3.18})| \leq \rho_{0}\|R\|_{\sigma}\|\mathcal{E}\|_{\sigma} .
$$

Combining the estimates for (3.15), (3.16), (3.17) and (3.18) we have

$$
\begin{equation*}
\left\langle\Lambda^{\sigma}\left(A_{0}\left(U^{\omega, n}\right) \mathcal{E}_{t}\right), \Lambda^{\sigma} \mathcal{E}\right\rangle \leq C n^{\sigma+1-s}\|\mathcal{E}\|_{H^{\sigma}}^{2}+C\|R\|_{\sigma}\|\mathcal{E}\|_{\sigma} \tag{3.22}
\end{equation*}
$$

where the constants depend upon $\rho_{0}, h_{0}, \gamma$ and $\sigma$.
We show that the residue $R$ satisfies the following estimate.
Proposition 2. If $\max \{1, s-2\}<\sigma<s-1$, then $\|R\|_{\sigma} \lesssim n^{\delta-2}$.

Proof. (3.6) the nonzero components of $R$ are

$$
\binom{R_{2}}{R_{3}}=\binom{\partial_{t} u^{\omega, n}+u^{\omega, n} \partial_{x} u^{\omega, n}+v^{\omega, n} \partial_{y} u^{\omega, n}}{\partial_{t} v^{\omega, n}+u \partial_{x} v^{\omega, n}+v^{\omega, n} \partial_{y} v^{\omega, n}} .
$$

3.5. Estimating $R_{2}$. We have (omitting the superscripts for brevity)

$$
\begin{aligned}
R_{2} & =u_{t}+u u_{x}+v u_{y}=\left(u_{1}+u_{2}\right)_{t}+\left(u_{1}+u_{2}\right)\left(u_{1}+u_{2}\right)_{x}+\left(v_{1}+v_{2}\right)\left(u_{1}+u_{2}\right)_{y} \\
& =\overbrace{u_{2, t}}^{0}+\overbrace{u_{1} u_{1, x}}^{1}+\overbrace{u_{1} u_{2, x}}^{2}+\overbrace{u_{2} u_{1, x}}^{3}+\overbrace{u_{2} u_{2, x}}^{4}+\overbrace{\overbrace{1} u_{1, y}}^{5}+\overbrace{v_{1} u_{2, y}}^{6}+\overbrace{v_{2} u_{1, y}}^{-7}+\overbrace{v_{2} u_{2, y}}^{8} .
\end{aligned}
$$

Now, three of these terms are zero by design, since $\operatorname{supp} u_{2}=\operatorname{supp} v_{2}=\operatorname{supp} S$ and $\phi_{2}^{\prime}=0=\phi_{2}^{\prime \prime}$ for $y \in \operatorname{supp} S$ :

$$
\begin{aligned}
& 2 \text { 2 } \equiv \frac{\omega}{n} \phi_{1}\left(\frac{x}{n^{\delta}}\right) \phi_{2}^{\prime}\left(\frac{y}{n^{\delta}}\right) \frac{1}{n^{\delta+s+1}} \partial_{x y} S=0 \\
& 3 \equiv u_{2} u_{1, x}=\frac{1}{n^{\delta+s+1}} \partial_{y} S \frac{\omega}{n^{1+\delta}} \phi_{1}^{\prime}\left(\frac{x}{n^{\delta}}\right) \phi_{2}^{\prime}\left(\frac{y}{n^{\delta}}\right)=0 \\
& 7 \equiv v_{2} u_{1, y}=-\frac{1}{n^{\delta+s+1}} \partial_{x} S\left(-\frac{\omega}{n^{1+\delta}}\right) \phi_{1}\left(\frac{x}{n^{\delta}}\right) \phi_{2}^{\prime \prime}\left(\frac{y}{n^{\delta}}\right)=0 \text {; }
\end{aligned}
$$

and another term takes a simpler form, since $\phi_{1}^{\prime} \equiv 1 \equiv \phi_{2}$ on the support of $S$ :

$$
6 \equiv v_{1} u_{2, y}=-\frac{\omega}{n} \phi_{1}^{\prime}\left(\frac{x}{n^{\delta}}\right) \phi_{2}\left(\frac{y}{n^{\delta}}\right) \frac{1}{n^{\delta+s+1}} \partial_{y} S=-\frac{\omega}{n^{\delta+s+2}} \partial_{y}^{2} S .
$$

From the form of the low-frequency and high-frequency terms, it is clear that differentiation of $u_{1}$ or $v_{1}$ with respect to either $x$ or $y$ improves the result by a factor of $n^{-\delta}$, as does differentiation of $S$ with respect to $x$; however, differentiation of $S$ with respect to $y$ introduces a term with an additional multiplicative factor of $n$. The amplitudes of the low- and high-frequency terms have been balanced so that the largest contributions due to this, in 0 and 6, cancel each other. This is exhibited in the proof of

Lemma 6 (Crucial cancellation). If $\varphi_{1}^{\prime} \equiv 1$ on $\operatorname{supp} \psi^{\prime}$ and $\varphi_{2} \equiv 1$ on $\operatorname{supp} \psi$, then

$$
\begin{align*}
u_{2, t}+v_{1} u_{2, y} & \equiv \frac{1}{n^{\delta+s+1}} \partial_{y}\left(\partial_{t}-\frac{\omega}{n} \partial_{y}\right) S \\
& =-\frac{\omega}{n^{2+2 \delta+s}} \psi\left(\frac{x}{n^{\delta}}\right)\left[\frac{1}{n^{\delta}} \psi^{\prime \prime}\left(\frac{y}{n^{\delta}}\right) \sin (n y+\omega t)+n \psi^{\prime}\left(\frac{y}{n^{\delta}}\right) \cos (n y+\omega t)\right], \tag{3.23}
\end{align*}
$$

and hence

$$
\|00+6\|_{\sigma} \leq C n^{\sigma-\delta-s-1} .
$$

Proof. Using $S(x, y, t)=\psi\left(n^{-\delta} x\right) \psi\left(n^{-\delta} y\right) \sin (n y+\omega t)$, we calculate

$$
\begin{aligned}
\left(\partial_{t}-\frac{\omega}{n} \partial_{y}\right) S & =\psi\left(\frac{x}{n^{\delta}}\right)\left(\partial_{t}-\frac{\omega}{n} \partial_{y}\right)\left(\psi\left(\frac{y}{n^{\delta}}\right) \sin (n y+\omega t)\right) \\
& =\psi\left(\frac{x}{n^{\delta}}\right)\left[\sin (n y+\omega t)\left(\partial_{t}-\frac{\omega}{n} \partial_{y}\right) \psi\left(\frac{y}{n^{\delta}}\right)+\psi\left(\frac{y}{n^{\delta}}\right)\left(\partial_{t}-\frac{\omega}{n} \partial_{y}\right) \sin (n y+\omega t)\right] \\
& =\psi\left(\frac{x}{n^{\delta}}\right)\left[\sin (n y+\omega t)\left(\partial_{t}-\frac{\omega}{n} \partial_{y}\right) \psi\left(\frac{y}{n^{\delta}}\right)\right] \\
& =-\frac{\omega}{n^{1+\delta}} \psi\left(\frac{x}{n^{\delta}}\right) \psi^{\prime}\left(\frac{y}{n^{\delta}}\right) \sin (n y+\omega t) .
\end{aligned}
$$

From this we obtain (3.23). Now it is a direct application of estimate (3.13) to complete the proof.

To complete the estimate for the $H^{\sigma}$ norm of $R_{2}$, we estimate the norms of $S$ and its derivatives. From

$$
S=\psi\left(\frac{x}{n^{\delta}}\right) \psi\left(\frac{y}{n^{\delta}}\right) \sin (n y+\omega t)
$$

and Lemma 4 we have $\|S\|_{\sigma} \lesssim n^{\sigma+\delta}$. Since differentiation with respect to $x$ scales the expression by $n^{-\delta}$ and differentiation with respect to $y$ scales it by $n$ (where we ignore the lower order contribution), we have

$$
\begin{equation*}
\left\|\partial_{x} S\right\|_{\sigma} \lesssim n^{\sigma}, \quad\left\|\partial_{y} S\right\|_{\sigma} \lesssim n^{\sigma+\delta+1}, \quad\left\|\partial_{x} \partial_{y} S\right\|_{\sigma} \lesssim n^{\sigma+1}, \quad\left\|\partial_{y}^{2} S\right\|_{\sigma} \lesssim n^{\sigma+\delta+2} . \tag{3.24}
\end{equation*}
$$

We also note the $H^{\sigma}$ bounds on $u_{1}$ and $v_{1}$ and their derivatives:

$$
\left\|u_{1}\right\|_{\sigma}=\left\|\frac{\omega}{n} \phi_{1}\left(\frac{x}{n^{\delta}}\right) \phi_{2}^{\prime}\left(\frac{x}{n^{\delta}}\right)\right\|_{\sigma} \lesssim n^{\delta-1}, \quad\left\|u_{1, x}\right\|_{\sigma} \lesssim \frac{1}{n},
$$

and the same bounds hold for $v_{1}$ and for the $y$ derivatives. With this we can find the remaining bounds for $R_{2}$ :

$$
\begin{aligned}
& \|\boxed{1}\|_{\sigma}=\left\|u_{1} u_{1, x}\right\|_{\sigma} \lesssim n^{\delta-2}, \\
& \|4\|_{\sigma}=\left\|u_{2} u_{2, x}\right\|_{\sigma}=\frac{1}{\left(n^{\delta+s+1}\right)^{2}}\left\|S_{y} S_{x y}\right\|_{\sigma} \lesssim n^{-2 s+2 \sigma-\delta}, \\
& \|5\|_{\sigma}=\left\|v_{1} u_{1, y}\right\|_{\sigma} \lesssim n^{\delta-2}, \\
& \|8\|_{\sigma}=\left\|v_{2} u_{2, y}\right\|_{\sigma}=\frac{1}{\left(n^{\delta+s+1}\right)^{2}}\left\|S_{x} S_{y y}\right\|_{\sigma} \lesssim n^{-2 s+2 \sigma-\delta} .
\end{aligned}
$$

Combining this with Lemma 6, we find the $H_{\sigma}$ norm of $R_{2}$ to be bounded by $n^{\alpha}$ where

$$
\alpha=\max \{\delta-2,-2(s-\sigma)-\delta, \sigma-\delta-s-1\} .
$$

Since $\sigma<s-1$, if we now choose $\delta \ll 1$, the largest exponent is $\delta-2$, so we have

$$
\begin{equation*}
\left\|R_{2}\right\|_{\sigma} \lesssim n^{\delta-2} . \tag{3.25}
\end{equation*}
$$

3.6. Estimating $R_{3}$. This goes the same way. (Again we omit the superscripts.)

$$
\begin{aligned}
R_{3} & =v_{t}+u v_{x}+v v_{y}=\left(v_{1}+v_{2}\right)_{t}+\left(u_{1}+u_{2}\right)\left(v_{1}+v_{2}\right)_{x}+\left(v_{1}+v_{2}\right)\left(v_{1}+v_{2}\right)_{y} \\
& =\overbrace{v_{2, t}}^{0}+\overbrace{u_{1} v_{1, x}}^{\boxed{1}}+\overbrace{u_{1} v_{2, x}}^{\boxed{2}}+\overbrace{u_{2} v_{1, x}}^{\boxed{3}}+\overbrace{u_{2} v_{2, x}}^{\boxed{4}}+\overbrace{v_{1} v_{1, y}}^{\boxed{5}}+\overbrace{v_{1} v_{2, y}}^{\boxed{6}}+\overbrace{v_{2} v_{1, y}}^{7}+\overbrace{v_{2} v_{2, y}}^{\boxed{8}} .
\end{aligned}
$$

Because $u_{1}, v_{1, x}$ and $v_{1, y}$ are zero on the support of $S$, we find that the terms 2,3 and 7 are again zero and (since $v_{1}$ is constant on $\operatorname{supp} S$ ), 6 reduces to $-\omega v_{2, y} / n$. This again gives us a cancellation between the highest order terms in 0 and 6 (we do not actually need it in the case of $R_{3}$ since the largest terms are already smaller by a factor of $n$ ). Specifically, using the identity in the proof of Lemma 6 .

$$
\begin{aligned}
0+6 & =v_{2, t}+v_{1} v_{2, y}=v_{2, t}-\frac{\omega}{n} v_{2, y}=-\frac{1}{n^{\delta+s+1}} \partial_{x}\left(\partial_{t}-\frac{\omega}{n} \partial_{y}\right) S \\
& =-\frac{1}{n^{\delta+s+1}} \partial_{x}\left(-\frac{\omega}{n^{1+\delta}} \psi\left(\frac{x}{n^{\delta}}\right) \psi^{\prime}\left(\frac{y}{n^{\delta}}\right) \sin (n y+\omega t)\right) \\
& =\frac{1}{n^{3 \delta+s+2}} \psi^{\prime}\left(\frac{x}{n^{\delta}}\right) \psi^{\prime}\left(\frac{y}{n^{\delta}}\right) \sin (n y+\omega t),
\end{aligned}
$$

and so

$$
\|\boxed{0}+6\| \|_{\sigma} \lesssim n^{-2 \delta-s+\sigma-2} .
$$

The estimates for the remaining terms are straightforward, as in the estimates for $R_{2}$. We use (3.24) and we need also $\left\|S_{x x}\right\|_{\sigma} \lesssim n^{\sigma-\delta}$ :

$$
\begin{aligned}
& \|\boxed{1}\|_{\sigma}=\left\|u_{1} v_{1, x}\right\|_{\sigma} \lesssim n^{\delta-2}, \\
& \|4\|_{\sigma}=\left\|u_{2} v_{2, x}\right\|_{\sigma}=\frac{1}{\left(n^{\delta+s+1}\right)^{2}}\left\|S_{y} S_{x x}\right\|_{\sigma} \lesssim n^{-2 s+2 \sigma-2 \delta-1}, \\
& \|5\|_{\sigma}=\left\|v_{1} v_{1, y}\right\|_{\sigma} \lesssim n^{\delta-2}, \\
& \|8\|\left\|_{\sigma}=\right\| v_{2} v_{2, y}\left\|_{\sigma}=\frac{1}{\left(n^{\delta+s+1}\right)^{2}}\right\| S_{x} S_{x y} \|_{\sigma} \lesssim n^{-2 s+2 \sigma-2 \delta-1} .
\end{aligned}
$$

Once again, the largest exponent is $\delta-2$, and so

$$
\begin{equation*}
\left\|R_{3}\right\|_{\sigma} \lesssim n^{\delta-2} . \tag{3.26}
\end{equation*}
$$

Combining estimates (3.25) and (3.26) completes the proof of Proposition 2.
To complete the proof of Theorem 6, first notice that from the definition $A_{0}\left(U^{\omega, n}\right) \geq c I$ for some positive constant $c$, and $A_{0}\left(U^{\omega, n}\right)$ is a constant matrix. Therefore, the $L^{2}$ inner product $\left\langle A_{0}\left(U^{\omega, n}\right) V, V\right\rangle$ defines an equivalent norm. Thus,

$$
\begin{equation*}
\frac{d}{d t}\|\mathcal{E}\|_{\sigma}^{2}=\frac{d}{d t}\left\langle\Lambda^{\sigma} \mathcal{E}, \Lambda^{\sigma} \mathcal{E}\right\rangle \approx \frac{d}{d t}\left\langle A_{0}\left(U^{\omega, n}\right) \Lambda^{\sigma} \mathcal{E}, \Lambda^{\sigma} \mathcal{E}\right\rangle \tag{3.27}
\end{equation*}
$$

Applying the derivative, we have

$$
\begin{align*}
\frac{d}{d t}\left\langle A_{0}\left(U^{\omega, n}\right) \Lambda^{\sigma} \mathcal{E}, \Lambda^{\sigma} \mathcal{E}\right\rangle & =2\left\langle A_{0}\left(U^{\omega, n}\right) \Lambda^{\sigma} \mathcal{E}_{t}, \Lambda^{\sigma} \mathcal{E}\right\rangle+\left\langle A_{0}\left(U^{\omega, n}\right)_{t} \Lambda^{\sigma} \mathcal{E}, \Lambda^{\sigma} \mathcal{E}\right\rangle  \tag{3.28}\\
& =2\left\langle\Lambda^{\sigma}\left(A_{0}\left(U^{\omega, n}\right) \mathcal{E}_{t}\right), \Lambda^{\sigma} \mathcal{E}\right\rangle
\end{align*}
$$

(since $A_{0}\left(U^{\omega, n}\right)$ is constant). This quantity was estimated in Section 3.4, substituting inequality (3.22) and applying Proposition 2 we have

$$
\begin{equation*}
2\|\mathcal{E}\|_{\sigma} \frac{d}{d t}\|\mathcal{E}\|_{\sigma} \approx \frac{d}{d t}\left\langle A_{0}\left(U^{\omega, n}\right) \Lambda^{\sigma} \mathcal{E}, \Lambda^{\sigma} \mathcal{E}\right\rangle \lesssim n^{\sigma+1-s}\|\mathcal{E}\|_{H^{\sigma}}^{2}+n^{\delta-2}\|\mathcal{E}\|_{\sigma} \tag{3.29}
\end{equation*}
$$

Dividing by $\|\mathcal{E}\|_{\sigma}$ in (3.29) gives the first inequality stated in the theorem. We apply Grönwall's inequality, [18, p 24]. The Grönwall estimate for

$$
z^{\prime}(t) \leq a z(t)+b, \quad z(0)=0,
$$

is

$$
z(t) \leq \frac{b}{a}\left(e^{a t}-1\right)
$$

Since here $a \simeq n^{\sigma-s+1}<C$, the upper bound for $e^{a t}$ is a constant that depends only on $T^{*}$, the time interval on which we are tracking the solution, and with $b \simeq n^{\delta-2}$, then $b / a$ gives the estimate in the theorem.

This completes the proof of the uniformity (in $n$ ) of the approximation of $U^{\omega, n}$ to the actual solution, $U_{\omega, n}$, for $\omega=1$ and $\omega=-1$. For $\varepsilon$ in equation (3.10) we have $3-(s-\sigma)-\delta>2$.
3.7. Nonuniform convergence. We are now prepared to complete the proof of nonuniform convergence, the final item, (44), in the program. We use a fact we proved in [28]: For a range of $\tau>s$, (specifically $s<\tau \leq\lfloor s\rfloor+1$, where $\lfloor\cdot\rfloor$ is the greatest integer function), the error in the $H^{\tau}$ norm is bounded by

$$
\begin{equation*}
\|\mathcal{E}\|_{\tau} \lesssim n^{\tau-s} . \tag{3.30}
\end{equation*}
$$

This uses the form of $\mathcal{E}$ and the bound in Lemma 4. Interpolation yields an estimate for the error in the $s$ norm.

Lemma 7. For any $s>2$ and $n \gg 1$, there exists an $\varepsilon>0$ such that

$$
\begin{equation*}
\|\mathcal{E}\|_{s} \lesssim n^{-\varepsilon} . \tag{3.31}
\end{equation*}
$$

Proof. From Proposition 1 we have

$$
\begin{equation*}
\|\mathcal{E}\|_{s} \leq\|\mathcal{E}\|_{\sigma}^{\alpha}\|\mathcal{E}\|_{\tau}^{\beta}, \quad \text { where } \quad \alpha=\frac{\tau-s}{\tau-\sigma}, \quad \beta=\frac{s-\sigma}{\tau-\sigma} . \tag{3.32}
\end{equation*}
$$

Using Theorem 6 and (3.30) we find

$$
\begin{equation*}
\|\mathcal{E}\|_{s} \lesssim\left(n^{\delta-3+s-\sigma}\right)^{\alpha}\left(n^{\tau-s}\right)^{\beta}=n^{(\tau-s)(\delta-3+2 s-2 \sigma) /(\tau-\sigma)} . \tag{3.33}
\end{equation*}
$$

By choosing $\max \{1, s-3 / 2+\delta\}<\sigma<s-1$, we obtain $\delta-3+2 s-2 \sigma<0$, which completes the proof.

We are now ready to give the proof of Theorem (1)
Proof. We estimate the difference between two actual solutions by the triangle inequality

$$
\begin{equation*}
\left\|U_{1, n}-U_{-1, n}\right\|_{s} \geq\left\|U^{1, n}-U^{-1, n}\right\|_{s}-\left\|U^{1, n}-U_{1, n}\right\|_{s}-\left\|U^{-1, n}-U_{-1, n}\right\|_{s} . \tag{3.34}
\end{equation*}
$$

From Lemma 7, the last two terms tend to zero as $n \rightarrow \infty$, and therefore, tracking the terms that do not tend to zero as $n \rightarrow \infty$,

$$
\begin{aligned}
\left\|U_{1, n}-U_{-1, n}\right\|_{s} & \geq \liminf _{n \rightarrow \infty}\left\|U^{1, n}-U^{-1, n}\right\|_{s} \\
& \geq \lim _{n \rightarrow \infty}\left\|n^{-\delta-s} \psi\left(n^{-\delta} x\right) \psi\left(n^{-\delta} y\right)(\cos (n y+t)-\cos (n y-t))\right\|_{s} \\
& =\lim _{n \rightarrow \infty}\left\|n^{-\delta-s} \psi\left(n^{-\delta} x\right) \psi\left(n^{-\delta} y\right) \cos (n y)\right\|_{s}|\sin (t)| \approx|\sin (t)| .
\end{aligned}
$$

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