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## LOCAL WEL LPOSEDNESS FOR

 THE 2+1-DIMENSIONAL MONOPOLE EQUATION
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#### Abstract

The space-time monopole equation on $\mathbb{R}^{2+1}$ can be derived by a dimensional reduction of the antiselfdual Yang-Mills equations on $\mathbb{R}^{2+2}$. It can be also viewed as the hyperbolic analog of Bogomolny equations. We uncover null forms in the nonlinearities and employ optimal bilinear estimates in the framework of wave-Sobolev spaces. As a result, we show the equation is locally wellposed in the Coulomb gauge for initial data sufficiently small in $H^{s}$ for $s>\frac{1}{4}$.


## 1. Introduction

In this paper we study local wellposedness of the Cauchy problem for the monopole equation on $\mathbb{R}^{2+1}$ Minkowski space in the Coulomb gauge. The space-time monopole equation can be derived by a dimensional reduction from the antiselfdual Yang-Mills equations on $\mathbb{R}^{2+2}$, and is given by

$$
\begin{equation*}
F_{A}=* D_{A} \phi, \tag{ME}
\end{equation*}
$$

where $F_{A}$ is the curvature of a one-form connection $A$ on $\mathbb{R}^{2+1}, D_{A} \phi$ is a covariant derivative of the Higgs field $\phi$, and $*$ is the Hodge star operator with respect to the Minkowski $\mathbb{R}^{2+1}$ metric. (ME) is a hyperbolic analog of Bogomolny equations, and was first introduced by Ward [1989] and discussed from the point of view of twistors. Ward [1999] also studied its soliton solutions. Recently, Dai, Terng and Uhlenbeck [2006] gave a broad survey on the space-time monopole equation. In particular, using the inverse scattering transform they have shown global existence and uniqueness up to a gauge transformation for small initial data in $W^{2,1}$. However, $L^{2}$ based wellposedness theory for this equation has not been investigated. The objective of this paper is to fill this gap by specifically treating the Cauchy problem for rough initial data in $H^{s}$.

Written in coordinates, (ME) is a system of first order hyperbolic partial differential equations. The unknowns are a pair $(A, \phi)$. If $(A, \phi)$ solve the equation, then so do

$$
\lambda A(\lambda t, \lambda x) \text { and } \lambda \phi(\lambda t, \lambda x), \quad \text { for any } \lambda>0 .
$$

This results in the critical exponent $s_{c}=0$. Since in general one expects local wellposedness for $s>s_{c}$ the goal would be to show (ME) is wellposed for $s>0$. Nevertheless, the two dimensions create an obstacle, which so far only allows $s>\frac{1}{4}$. We explain this now. In Section 4 we choose a Coulomb gauge, and

[^0]reformulate (ME) as a system of semilinear wave equations coupled with an elliptic equation, to which we refer as auxiliary monopole equations (aME). Schematically it looks as follows:
\[

$$
\begin{align*}
\square u & =\mathscr{B}_{+}\left(\partial u, \partial v, A_{0}\right), \\
\square v & =\mathscr{B}_{-}\left(\partial u, \partial v, A_{0}\right),  \tag{aME}\\
\triangle A_{0} & =\mathscr{C}\left(\partial u, \partial v, A_{0}\right),
\end{align*}
$$
\]

where $\mathscr{B}_{ \pm}, \mathscr{C}$ are bilinear forms, ${ }^{1} A_{0}$ is the temporal part of the connection $A, \partial u, \partial v$ denote space-time derivatives of $u$ and $v$ respectively, and are given in terms of $\phi$ and the spatial part of $A$. As a result, showing wellposedness of (ME) for $s>0$ can follow from showing (aME) is wellposed for $s>1$ (see Theorem 4.1). Also, the most difficult nonlinearity that we have to handle is contained in $\mathscr{B}_{ \pm}\left(\partial u, \partial v, A_{0}\right)$. Luckily, it exhibits a structure of a null form. There are two standard null forms:

$$
Q_{0}(u, v)=-\partial_{t} u \partial_{t} v+\nabla u \cdot \nabla v, \quad Q_{\alpha \beta}(u, v)=\partial_{\alpha} u \partial_{\beta} v-\partial_{\beta} u \partial_{\alpha} v .
$$

The null condition was introduced by Klainerman [1984], and it was first applied to produce better local wellposedness results for wave equations with a null form by Klainerman and Machedon [1993]. Indeed, in low dimensions, for these kind of nonlinearities one can assume much less regularity of the initial data than for the general products. Counterexamples for general products were shown by Lindblad [1996]. We uncover the null form $Q_{\alpha \beta}$ in our system of wave equations as well as a new type of a null form which is related to $Q_{\alpha \beta}$. Unfortunately, the results in two spatial dimensions for $Q_{\alpha \beta}$ are not as optimal as they are in higher dimensions or as they are for $Q_{0}$. In fact, the best result in literature so far for $Q_{\alpha \beta}$ in two dimensions is due to Zhou [1997]. He establishes local wellposedness for initial data in $H^{s} \times H^{s-1}$ for $s>\frac{5}{4}$. In addition, by examining the first iterate Zhou shows that this is as close as one can get to the critical level using iteration methods. ${ }^{2}$ On the other hand, for dimensions $n \geq 3$ Klainerman and Machedon [1996] showed almost optimal local wellposedness in $H^{s} \times H^{s-1}$ for $s>n / 2$. The articles [Klainerman and Machedon 1995; Klainerman and Selberg 2002] give equally satisfying results for $Q_{0}$, and in all dimensions $n \geq 2$.

Now, one of the nonlinearities in the system (aME) is $Q_{\alpha \beta}$, so showing (aME) is locally wellposed for $s>\frac{5}{4}$ would be sharp by iteration methods. This is what we do, and as a result we obtain local wellposedness of (ME) in the Coulomb gauge for $s>\frac{1}{4}$ (see the Main Theorem below). However, (aME) is not exactly (ME), so we hope to treat (ME) directly in the near future and improve the results. What should be mentioned here is that we have considered other traditional gauges such as Lorentz and Temporal, but they have not been as nearly useful as the Coulomb gauge. Perhaps other, less traditional gauges could be used. Moreover, we note that even the estimates involving the temporal variable $A_{0}$ seem to require $s>\frac{1}{4}$.

We include a brief discussion about global wellposedness. As already mentioned, in [Dai et al. 2006] the inverse scattering transform is used to show global existence and uniqueness up to a gauge transformation for small data in $W^{2,1}$. To extend it to global wellposedness in $L^{2}$ based theory, we would like to benefit from the local result in this article. It is not immediately clear how this can be accomplished

[^1]since, for example, the energy functional is not positive definite and, in fact, it vanishes for the solutions of (ME) [Ward 1989]. Global wellposedness is an interesting question, and we would like to investigate it in the future.

The main result of this paper is contained in the following theorem.
Main Theorem. Let $\frac{1}{4}<s<\frac{1}{2}$ and $r \in(0,2 s]$ and consider the space-time monopole equation

$$
\begin{equation*}
F_{A}=* D_{A} \phi, \tag{ME}
\end{equation*}
$$

with initial data

$$
\left.\left(A_{1}, A_{2}, \phi\right)\right|_{t=0}=\left(a_{1}, a_{2}, \phi_{0}\right)
$$

then (ME) in the Coulomb gauge is locally wellposed for initial data sufficiently small in $H^{s}\left(\mathbb{R}^{2}\right)$ in the following sense.

- Local existence: For all $a_{1}, a_{2}, \phi_{0} \in H^{s}\left(\mathbb{R}^{2}\right)$ sufficiently small there exist $T>0$ depending continuously on the norm of the initial data, and functions

$$
A_{0} \in C_{b}\left([0, T], \dot{H}^{r}\right), \quad A_{1}, A_{2}, \phi \in C_{b}\left([0, T], H^{s}\right)
$$

which solve (ME) in the Coulomb gauge on $[0, T] \times \mathbb{R}^{2}$ in the sense of distributions and such that the initial conditions are satisfied.

- Uniqueness: If $T>0$ and $(A, \phi)$ and $\left(A^{\prime}, \phi^{\prime}\right)$ are two solutions of (ME) in the Coulomb gauge on $(0, T) \times \mathbb{R}^{2}$ belonging to

$$
C_{b}\left([0, T], \dot{H}^{r}\right) \times\left(H_{T}^{s, \theta}\right)^{3}
$$

with the same initial data, then $(A, \phi)=\left(A^{\prime}, \phi^{\prime}\right)$ on $(0, T) \times \mathbb{R}^{2}$.

- Continuous dependence on the initial data: For any $a_{1}, a_{2}, \phi_{0} \in H^{s}\left(\mathbb{R}^{2}\right)$ there is a neighborhood $U$ of $a_{1}, a_{2}, \phi_{0}$ in $\left(H^{s}\left(\mathbb{R}^{2}\right)\right)^{3}$ such that the solution map $\left(a, \phi_{0}\right) \rightarrow(A, \phi)$ is continuous from $U$ into $C_{b}\left([0, T], \dot{H}^{r}\right) \times\left(C_{b}\left([0, T], H^{s}\right)\right)^{3}$.
Remark 1.1. The spaces $H_{T}^{s, \theta}$ are defined in Section 2B.
Remark 1.2. The initial data does not have to be given in the Coulomb gauge. See Theorem 3.3.
Remark 1.3. There are two reasons for the requirement of the small initial data. First, the construction of the global Coulomb gauge requires an assumption on the size of the data (see Section 3B). The second obstacle comes from the elliptic equation for $A_{0}$ in (aME), and including $A_{0}$ in the Picard iteration. See Remark 4.2 for further discussion.
Remark 1.4. We do not prescribe initial data for $A_{0}$, because when $A$ is in the Coulomb gauge, $A_{0}(t)$ can be determined at any time by solving the elliptic equation. See Section 4 for more details.
Remark 1.5. To simplify the exposition, in this paper we assume $\frac{1}{4}<s<\frac{1}{2}$. See [Czubak 2008] for all $s>\frac{1}{4}$. In general, the higher the value of $s$, the less delicate the estimates have to be. We have a uniform way to handle all $s>\frac{1}{4}$ for the estimates involving the null forms (see Section 5B1 for a discussion). Therefore the reason for restricting the range of $s$ is rather due to the technicalities of the estimates for $A_{0}$ (Theorems 5.3 and 5.5 and Corollaries 5.4 and 5.6) and the regularity of the gauge transformations (Lemma 3.1). The technicalities are not very interesting and are handled in [Czubak 2008] with similar arguments as those presented here.

The outline of the paper is as follows. Section 2 sets notation, introduces spaces, and estimates used throughout the paper. In Section 3 we take a closer look at the equations and discuss gauge transformations. In Section 4 we rewrite (ME) as a system of wave equations coupled with an elliptic equation. We also show local wellposedness of the new system implies local wellposedness of (ME) in the Coulomb gauge. Section 5 is devoted to the proof of the Main Theorem, which is reduced to establishing estimates (5-4)-(5-8).

## 2. Preliminaries

First we establish notation, then we introduce function spaces as well as estimates used.
2A. Notation. The expression $a \lesssim b$ means $a \leq C b$ for some positive constant $C$. A point in the $2+1$ dimensional Minkowski space is written as $(t, x)=\left(x^{\alpha}\right)_{0 \leq \alpha \leq 2}$. Greek indices range from 0 to 2 , and Roman indices range from 1 to 2 . We raise and lower indices with the Minkowski metric diag $(-1,1,1)$. We write $\partial_{\alpha}=\partial_{x^{\alpha}}$ and $\partial_{t}=\partial_{0}$, and we also use the Einstein notation. Therefore, $\partial^{i} \partial_{i}=\Delta$, and $\partial^{\alpha} \partial_{\alpha}=$ $-\partial_{t}^{2}+\Delta=\square$. When we refer to spatial and time derivatives of a function $f$, we write $\partial f$, and when we consider only spatial derivatives of $f$, we write $\nabla f$. Finally, $d$ denotes the exterior differentiation operator and $d^{*}$ its dual given by $d^{*}=(-1)^{k} * * * d *$, where $*$ is the Hodge $*$ operator (see, for example, [Roe 1998]) and $k$ comes from $d^{*}$ acting on some given $k$-form. It will be clear from the context, when * and $d^{*}$ operators act with respect to the Minkowski metric and when with respect to the Euclidean metric. For the convenience of the reader we include the following: with respect to the Euclidean metric on $\mathbb{R}^{2}$ we have

$$
* d x=d y, \quad * d y=-d x, \quad * 1=d x \wedge d y
$$

and with respect to the $\operatorname{diag}(-1,1,1)$ metric on $\mathbb{R}^{2+1}$,

$$
* d t=d x \wedge d y, \quad * d x=d t \wedge d y, \quad * d y=-d t \wedge d x
$$

2B. Function spaces. We use Picard iteration. Here we introduce the spaces, in which we are going to perform the iteration ${ }^{3}$. First we define the following Fourier multiplier operators

$$
\begin{gathered}
\widehat{\Lambda^{\alpha}} f(\xi)=\left(1+|\xi|^{2}\right)^{\alpha / 2} \hat{f}(\xi), \quad \widehat{\Lambda_{+}^{\alpha} u}(\tau, \xi)=\left(1+\tau^{2}+|\xi|^{2}\right)^{\alpha / 2} \hat{u}(\tau, \xi), \\
\widehat{\Lambda_{-}^{\alpha}} u(\tau, \xi)=\left(1+\frac{\left(\tau^{2}-|\xi|^{2}\right)^{2}}{1+\tau^{2}+|\xi|^{2}}\right)^{\alpha / 2} \hat{u}(\tau, \xi),
\end{gathered}
$$

where the symbol of $\Lambda_{-}^{\alpha}$ is comparable to $(1+||\tau|-|\xi||)^{\alpha}$. The corresponding homogeneous operators are denoted by $D^{\alpha}, D_{+}^{\alpha}, D_{-}^{\alpha}$, respectively.

Now, the spaces of interest are the wave-Sobolev spaces, $H^{s, \theta}$ and $\mathscr{H}^{s, \theta}$, given by ${ }^{4}$

$$
\|u\|_{H^{s, \theta}}=\left\|\Lambda^{s} \Lambda_{-}^{\theta} u\right\|_{L^{2}\left(\mathbb{R}^{2+1}\right)}, \quad\|u\|_{\mathscr{H}, \theta}=\|u\|_{H^{s, \theta}}+\left\|\partial_{t} u\right\|_{H^{s-1, \theta}} .
$$

[^2]An equivalent norm for $\mathscr{L}^{s, \theta}$ is $\|u\|_{\mathscr{H} s, \theta}=\left\|\Lambda^{s-1} \Lambda_{+} \Lambda_{-}^{\theta} u\right\|_{L^{2}\left(\mathbb{R}^{2+1}\right)}$. By results in [Selberg 1999] if $\theta>1 / 2$, we have

$$
\begin{align*}
& H^{s, \theta} \hookrightarrow C_{b}\left(\mathbb{R}, H^{s}\right),  \tag{2-1}\\
& \mathscr{H}^{s, \theta} \hookrightarrow C_{b}\left(\mathbb{R}, H^{s}\right) \cap C_{b}^{1}\left(\mathbb{R}, H^{s-1}\right) . \tag{2-2}
\end{align*}
$$

This is a crucial fact needed to localize our solutions in time. We denote the restrictions to the time interval $[0, T]$ by

$$
H_{T}^{s, \theta} \quad \text { and } \quad \mathscr{H}_{T}^{s, \theta},
$$

respectively.
2C. Estimates used. Throughout the paper we use the following estimates:

$$
\begin{equation*}
\left\|D^{-\sigma}(u v)\right\|_{L_{t}^{p} L_{x}^{q}} \lesssim\|u\|_{H^{s, \theta}}\|v\|_{H^{s, \theta}} \tag{2-3}
\end{equation*}
$$

is a theorem established by Klainerman and Tataru [1999] for the space-time operator $D_{+}$. The proof for the spatial operator $D$ is in [Selberg 1999]. There are several conditions that $\sigma, p, q$ have to satisfy, and they are listed in Section 5D, where we discuss the application of the estimate. Further,

$$
\begin{align*}
\|u\|_{L_{t}^{p} L_{x}^{2}} \lesssim\|u\|_{H^{0, \theta}} \leq \infty & & \text { if } 2 \leq p \leq \infty, \theta>\frac{1}{2},  \tag{2-4}\\
\|u\|_{L_{t}^{p} L_{x}^{q}} \lesssim\|u\|_{H^{1-2 / q-1 / p, \theta}} & & \text { if } 2 \leq p \leq \infty, 2 \leq q<\infty, 2 / p \leq \frac{1}{2}-1 / q, \theta>\frac{1}{2},  \tag{2-5}\\
\|u v\|_{L_{t, x}^{2}} \lesssim\|u\|_{H^{a, \alpha}}\|v\|_{H^{b, \beta}} & & \text { if } a, b, \alpha, \beta \geq 0, a+b>1, \alpha+\beta>\frac{1}{2} . \tag{2-6}
\end{align*}
$$

Estimate (2-4) can be proved by interpolation between $H^{0, \theta} \hookrightarrow L_{t, x}^{2}$ and (2-1) with $s=0$. Estimate (2-5) is a two-dimensional case of Theorem D in [Klainerman and Selberg 2002]. Finally, (2-6) is a special case of the proposition in [Klainerman and Selberg 2002, Appendix A.2].

## 3. A closer look at the monopole equations

3A. Derivation and background. Electric charge is quantized, which means that it appears in integer multiples of an electron. This is called the principle of quantization and has been observed in nature. The only theoretical proof so far was presented by Paul Dirac [1931]. In the proof Dirac introduced the concept of a magnetic monopole, of an isolated point-source of a magnetic charge. Despite extensive research, magnetic monopoles have not been (yet) found in nature. We refer to magnetic monopoles as Euclidean monopoles. The Euclidean monopole equation has exactly the same form as our space-time monopole equation (ME),

$$
F_{A}=* D_{A} \phi,
$$

with the exception that $*$ acts here with respect to the Euclidean metric and the base manifold is $\mathbb{R}^{3}$ instead of $\mathbb{R}^{2+1}$. The Euclidean monopole equations are also referred to as Bogomolny equations. For more on Euclidean monopoles we refer the reader to [Jaffe and Taubes 1980] and [Atiyah and Hitchin 1988]. In this paper we study the space-time monopole equation, which was first introduced by Ward [1989]. Both the Euclidean and the space-time monopole equations are examples of integrable systems
and have an equivalent formulation as a Lax pair. This and much more can be found in [Dai et al. 2006].
Given a space-time monopole equation

$$
\begin{equation*}
F_{A}=* D_{A} \phi \tag{ME}
\end{equation*}
$$

the unknowns are a pair $(A, \phi) . A$ is a connection 1-form given by

$$
\begin{equation*}
A=A_{0} d t+A_{1} d x+A_{2} d y, \quad \text { where } A_{\alpha}: \mathbb{R}^{2+1} \rightarrow \mathfrak{g} \tag{3-1}
\end{equation*}
$$

Here $\mathfrak{g}$ is the Lie algebra of a Lie group $G$, which is typically taken to be a matrix group $S U(n)$ or $U(n)$. In this paper we consider $G=S U(n)$, but everything we say here should generalize to any compact Lie group.

To be more general we could say $A$ is a connection on a principal $G$-bundle. Then observe that the $G$-bundle we deal here with is a trivial bundle $\mathbb{R}^{2+1} \times G$.

Next, $\phi$ is a section of a vector bundle associated to the $G$-bundle by a representation. We use the adjoint representation. Since we have a trivial bundle, we can just think of the Higgs field $\phi$ as a map from $\mathbb{R}^{2,1} \rightarrow \mathfrak{g}$.
$F_{A}$ is the curvature of $A$. It is a Lie algebra valued 2-form on $\mathbb{R}^{2+1}$

$$
\begin{equation*}
F_{A}=\frac{1}{2} F_{\alpha \beta} d x^{\alpha} \wedge d x^{\beta}, \quad \text { where } F_{\alpha \beta}=\partial_{\alpha} A_{\beta}-\partial_{\beta} A_{\alpha}+\left[A_{\alpha}, A_{\beta}\right] \tag{3-2}
\end{equation*}
$$

where $[\cdot, \cdot]$ denotes the Lie bracket, which for matrices can be thought of simply as $[X, Y]=X Y-Y X$. When we write $[\phi, B]$, where $B$ is a 1 -form, we mean

$$
\begin{equation*}
[\phi, B]=\left[\phi, B_{i}\right] d x^{i} \quad \text { and } \quad[B, C]=\frac{1}{2}\left[B_{i}, C_{j}\right] d x^{i} \wedge d x^{j}, \text { for two 1-forms } B, C . \tag{3-3}
\end{equation*}
$$

In the physics language, frequently adopted by the mathematicians, $A$ is called a gauge potential, $\phi$ a scalar field and $F_{A}$ is called an electromagnetic field.

Next, $D_{A}$ is the covariant exterior derivative associated to $A$, and $D_{A} \phi$ is given by

$$
\begin{equation*}
D_{A} \phi=D_{\alpha} \phi d x^{\alpha}, \quad \text { where } D_{\alpha} \phi=\partial_{\alpha} \phi+\left[A_{\alpha}, \phi\right] . \tag{3-4}
\end{equation*}
$$

The space-time monopole equation (ME) is obtained by a dimensional reduction of the antiselfdual Yang-Mills equations on $\mathbb{R}^{2+2}$, given by

$$
\begin{equation*}
F_{A}=-* F_{A} \tag{ASDYM}
\end{equation*}
$$

If the curvature of a connection $A$ satisfies (ASDYM), then $A$ is called an antiselfdual connection. (The corresponding selfdual Yang-Mills equation is

$$
\begin{equation*}
F_{A}=* F_{A} ; \tag{SDYM}
\end{equation*}
$$

in either case $F_{A}$ satisfies the Yang-Mills equation $D_{A} * F=0$, since $F_{A}= \pm * F_{A}$ implies

$$
D_{A} * F= \pm D_{A} F=0,
$$

as can be seen from the second Bianchi identity, or by direct computation from (3-2).)
Both (ASDYM) and (SDYM) are known to give rise to many different integrable equations (see [Ward 1985; Ablowitz et al. 2003] and references therein). In particular, the $2+2$ signature is used to derive both KDV and NLS in [Mason and Sparling 1989], and harmonic maps from $\mathbb{R}^{2}$ and $\mathbb{R}^{1+1}$ into a Lie
group in [Uhlenbeck 1992]. Also see [Ward 1989], where Einstein's vacuum equation for cylindrically symmetric space-times, and the sine-Gordon equation are derived from (ME).

We now present the details of the derivation of (ME) from (ASDYM) outlined in [Dai et al. 2006]. Let

$$
d x_{1}^{2}+d x_{2}^{2}-d x_{3}^{2}-d x_{4}^{2}
$$

be a metric on $\mathbb{R}^{2+2}$, then in coordinates (ASDYM) is

$$
\begin{equation*}
F_{12}=-F_{34}, \quad F_{13}=-F_{24}, \quad F_{23}=F_{14} \tag{3-5}
\end{equation*}
$$

The next step is the dimensional reduction, where we assume the connection $A$ is independent of $x_{3}$, and we let $A_{3}=\phi$. Then (3-5) becomes

$$
\begin{equation*}
D_{0} \phi=F_{12}, \quad D_{1} \phi=F_{02}, \quad D_{2} \phi=F_{10}, \tag{3-6}
\end{equation*}
$$

where we use index 0 instead of 4 . This is exactly (ME) written out in components.
Remark 3.1. Equivalently we could write (ME) as

$$
\begin{equation*}
F_{\alpha \beta}=-\varepsilon_{\alpha \beta \gamma} D^{\gamma} \phi, \tag{3-7}
\end{equation*}
$$

where $\varepsilon_{\alpha \beta \gamma}$ is a completely antisymmetric tensor with $\varepsilon_{012}=1$, and where we raise the index $\gamma$ using the Minkowski metric. We choose to work with the Hodge operator $*$ as it simplifies our task in Section 4.

Following [Dai et al. 2006], there is another way to write (ME), which is very useful for computations. (ME) is an equation involving 2 -forms on both sides. By taking the parts corresponding to $d t \wedge d x$ and $d t \wedge d y$ on the one hand, and the parts corresponding to $d x \wedge d y$ on the other, we obtain, respectively,

$$
\begin{align*}
\partial_{t} A+\left[A_{0}, A\right]-d A_{0} & =* d \phi+[* A, \phi],  \tag{3-8}\\
d A+[A, A] & =*\left(\partial_{t} \phi+\left[A_{0}, \phi\right]\right) . \tag{3-9}
\end{align*}
$$

Observe that now operators $d$ and $*$ act only with respect to the spatial variables. Similarly, $A$ now denotes only the spatial part of the connection, i.e., $A=\left(A_{1}, A_{2}\right)$. Moreover, (3-8) is an equation involving 1 -forms, and (3-9) involves 2 -forms.

3B. Gauge transformations. (ME) is invariant under gauge transformations. Indeed, if we have a smooth map $g$, with compact support such that $g: \mathbb{R}^{2+1} \rightarrow G$, and

$$
\begin{equation*}
A \rightarrow A_{g}=g A g^{-1}+g d g^{-1}, \quad \phi \rightarrow \phi_{g}=g \phi g^{-1} \tag{3-10}
\end{equation*}
$$

then a computation shows $F_{A} \rightarrow g F_{A} g^{-1}$ and $D_{A} \phi \rightarrow g D_{A} \phi g^{-1}$. Therefore if a pair $(A, \phi)$ solves (ME), so does $\left(A_{g}, \phi_{g}\right)$.

We would like to discuss the regularity of the gauge transformations. If $A \in X, \phi \in Y$ where $X, Y$ are some Banach spaces, the smoothness and compact support assumption on $g$ can be lowered just enough so the gauge transformation defined above is a continuous map from $X$ back into $X$, and from $Y$ back into $Y$. First note that since we are mapping into a compact Lie group, we can assume $g \in L_{t, x}^{\infty}$ and $\|g\|_{L_{t, x}^{\infty}}=\left\|g^{-1}\right\|_{L_{t, x}^{\infty}}$. Next, note that the Main Theorem produces a solution so that $\phi$ and the spatial parts of the connection $A_{1}, A_{2} \in C_{b}\left(I, H^{s}\right), \frac{1}{4}<s<\frac{1}{2}$, and $A_{0} \in C_{b}\left(I, \dot{H}^{r}\right), r \in(0,2 s]$.

Lemma 3.1. Let $0<\alpha<1$, and $Y=C_{b}\left(I, \dot{H}^{1} \cap \dot{H}^{\alpha+1}\right) \cap L^{\infty}$, then the gauge action is a continuous map from

$$
\begin{equation*}
C_{b}\left(I, H^{\alpha}\right) \times Y \rightarrow C_{b}\left(I, H^{\alpha}\right), \quad(h, g) \mapsto g h g^{-1}+g d g^{-1} \tag{3-11}
\end{equation*}
$$

and the following estimate holds:

$$
\begin{equation*}
\left\|h_{g}\right\|_{C_{b}\left(I, H^{\alpha}\right)} \lesssim\left(\|h\|_{C_{b}\left(I, H^{\alpha}\right)}+1\right)\|g\|_{Y}^{2} \tag{3-12}
\end{equation*}
$$

Proof. The continuity of the map easily follows from the inequalities we obtain below. Next, for fixed $t$ we have

$$
\left\|g(t) h(t) g^{-1}(t)+g(t) d g^{-1}(t)\right\|_{H^{\alpha}} \lesssim\left\|g h g^{-1}\right\|_{L^{2}}+\left\|D^{\alpha}\left(g h g^{-1}\right)\right\|_{L^{2}}+\left\|g d g^{-1}\right\|_{H^{\alpha}}
$$

where for the ease of notation we eliminated writing of the variable $t$ on the right side of the inequality. The first term is bounded by $\|h(t)\|_{H^{\alpha}}\|g\|_{L^{\infty}}^{2}$. For the second one we have

$$
\left\|D^{\alpha}\left(g h g^{-1}\right)\right\|_{L^{2}} \lesssim\left\|D^{\alpha} g h\right\|_{L^{2}}\|g\|_{L^{\infty}}+\left\|h D^{\alpha} g^{-1}\right\|_{L^{2}}\|g\|_{L^{\infty}}+\|h\|_{\dot{H}^{\alpha}}\|g\|_{L^{\infty}}^{2}
$$

It is enough to only look at the first term since $g$ and $g^{-1}$ have the same regularity. By Hölder's inequality and Sobolev embedding

$$
\begin{equation*}
\left\|D^{\alpha} g h\right\|_{L^{2}} \leq\left\|D^{\alpha} g\right\|_{L^{2 / \alpha}}\|h\|_{L^{(1 / 2-\alpha / 2)^{-1}}} \lesssim\|g\|_{\dot{H}^{1}}\|h\|_{\dot{H}^{\alpha}} \tag{3-13}
\end{equation*}
$$

where we use that $\frac{\alpha}{2}=\frac{1}{2}-\frac{1-\alpha}{2}$. Finally for the last term we have

$$
\begin{equation*}
\left\|g d g^{-1}\right\|_{H^{\alpha}} \lesssim\|g\|_{\dot{H}^{1}}\|g\|_{L^{\infty}}+\left\|D^{\alpha} g d g^{-1}\right\|_{L^{2}}+\|g\|_{\dot{H}^{\alpha+1}}\|g\|_{L^{\infty}} \tag{3-14}
\end{equation*}
$$

and we are done if we observe that the second term can be handled exactly as in (3-13).
Remark 3.2. We assume $0<\alpha<1$ since this is the case we need. However it is not difficult to see the lemma still holds with $\alpha=0$ or $\alpha \geq 1$ [Czubak 2008].

From the lemma, we trivially obtain the following corollary.
Corollary 3.2. Let $0<r, s<1$,

$$
X=C_{b}\left(I, \dot{H}^{r}\right) \times C_{b}\left(I, H^{s}\right) \times C_{b}\left(I, H^{s}\right), \quad \text { and } \quad Y=C_{b}\left(I, \dot{H}^{1} \cap \dot{H}^{s+1} \cap \dot{H}^{r+1}\right) \cap L^{\infty} .
$$

Then the gauge action is a continuous map from

$$
\begin{equation*}
X \times Y \rightarrow X, \quad\left(A_{0}, A_{1}, A_{2}\right) \mapsto A_{g} \tag{3-15}
\end{equation*}
$$

as well as from

$$
\begin{equation*}
C_{b}\left(I, H^{s}\right) \times Y \rightarrow C_{b}\left(I, H^{s}\right), \quad \phi \mapsto \phi_{g}=g \phi g^{-1} \tag{3-16}
\end{equation*}
$$

and the following estimates hold:

$$
\begin{align*}
\left\|A_{g}\right\|_{X} & \lesssim\left(1+\|A\|_{X}\right)\|g\|_{Y}^{2}  \tag{3-17}\\
\left\|\phi_{g}\right\|_{C_{b}\left(I, H^{s}\right)} & \lesssim\|\phi\|_{C_{b}\left(I, H^{s}\right)}\|g\|_{Y}^{2} .
\end{align*}
$$

Since in this paper we work in the Coulomb gauge, we ask: given any initial data $a_{1}, a_{2}, \phi_{0} \in H^{s}\left(\mathbb{R}^{2}\right)$, can we find a gauge transformation so that the initial data is placed in the Coulomb gauge? Dell'Antonio and Zwanziger [1991] produce a global $\dot{H}^{1}$ Coulomb gauge using variational methods. Here, we also require $g \in \dot{H}^{s+1}$, and two dimensions are tricky. Fortunately, if the initial data is small, we can obtain a global gauge with the additional regularity as needed. This has been studied by the author and Uhlenbeck for two dimensions and higher; the result in two dimensions is the following:
Theorem 3.3 [Czubak and Uhlenbeck $\geq 2010$ ]. Let $s>0$. Given $A(0)=a$ sufficiently small in

$$
H^{s}\left(\mathbb{R}^{2}\right) \times H^{s}\left(\mathbb{R}^{2}\right)
$$

there exists a gauge transformation $g \in \dot{H}^{s+1}\left(\mathbb{R}^{2}\right) \cap \dot{H}^{1}\left(\mathbb{R}^{2}\right) \cap L^{\infty}$ such that $\partial^{i}\left(g a_{i} g^{-1}+g \partial_{i} g^{-1}\right)=0$.

## 4. The monopole equation in the Coulomb gauge as a system of wave and elliptic equations

We begin by rewriting the monopole equation in the Coulomb gauge as a system of wave equations coupled with an elliptic equation. We refer to the new system as the auxiliary monopole equations (aME). Then we establish that local wellposedness (LWP) for (ME) in the Coulomb gauge can be obtained from LWP of (aME).

4A. Derivation of (aME) from (ME). Suppose we have initial data

$$
\begin{equation*}
\left.A_{i}\right|_{t=0}=a_{i} \quad \text { for } i=1,2 \quad \text { and }\left.\quad \phi\right|_{t=0}=\phi_{0}, \tag{4-1}
\end{equation*}
$$

where $\partial^{i} a_{i}=0$. Recall equations (3-8) and (3-9):

$$
\begin{align*}
\partial_{t} A+\left[A_{0}, A\right]-d A_{0} & =* d \phi+[* A, \phi]  \tag{4-2}\\
d A+[A, A] & =*\left(\partial_{t} \phi+\left[A_{0}, \phi\right]\right) \tag{4-3}
\end{align*}
$$

where $d$ and $*$ act only with respect to the spatial variables, and $A$ denotes only the spatial part of the connection. If we impose the Coulomb gauge condition, then

$$
\begin{equation*}
d^{*} A=0 \tag{4-4}
\end{equation*}
$$

By equivalence of closed and exact forms on $\mathbb{R}^{n}$, we can further suppose that

$$
\begin{equation*}
A=* d f \tag{4-5}
\end{equation*}
$$

for some $f: \mathbb{R}^{2+1} \rightarrow \mathfrak{g}$. Observe that

$$
\begin{equation*}
d * d f=\Delta f d x \wedge d y, \quad[* d f, * d f]=[d f, d f]=1 / 2\left[\partial_{i} f, \partial_{j} f\right] d x^{i} \wedge d x^{j} \tag{4-6}
\end{equation*}
$$

and $* * \omega=-\omega$ for a one-form on $\mathbb{R}^{2}$. Hence (4-2) and (4-3) become

$$
\begin{align*}
\partial_{t} * d f+\left[A_{0}, * d f\right]-d A_{0} & =* d \phi-[d f, \phi],  \tag{4-7}\\
\Delta f+\left[\partial_{1} f, \partial_{2} f\right] & =\partial_{t} \phi+\left[A_{0}, \phi\right] . \tag{4-8}
\end{align*}
$$

Take $d^{*}$ of (4-7) to obtain

$$
\Delta A_{0}=-d^{*}\left[A_{0}, * d f\right]-d^{*}[d f, \phi]
$$

This is the elliptic equation in (aME). To uncover the wave equations we proceed as follows. Take $d$ of (4-7)

$$
\begin{equation*}
\partial_{t} \Delta f+\partial^{j}\left[A_{0}, \partial_{j} f\right]=\Delta \phi+\partial_{2}\left[\partial_{1} f, \phi\right]-\partial_{1}\left[\partial_{2} f, \phi\right] . \tag{4-9}
\end{equation*}
$$

Consider (4-9) and (4-8) on the spatial Fourier transform side:

$$
\begin{align*}
-\partial_{t}|\xi|^{2} \hat{f}+|\xi|^{2} \hat{\phi} & =i\left(\xi_{2}\left[\widehat{\partial_{1} f, \phi}\right]-\xi_{1}\left[\widehat{\partial_{2} f, \phi}\right]-\xi_{j}\left[\widehat{A_{0}, \partial_{j}} f\right]\right),  \tag{4-10}\\
-|\xi|^{2} \hat{f}-\partial_{t} \hat{\phi} & =-\left[\widehat{\partial_{1} f, \partial_{2}} f\right]+\left[\widehat{A_{0}, \phi}\right] . \tag{4-11}
\end{align*}
$$

This allows us to write (4-10) and (4-11) as a system for $\phi$ and $d f$ :

$$
\begin{align*}
\left(\partial_{t}-i|\xi|\right)(\hat{\phi}+i|\xi| \hat{f}) & =-\hat{\mathscr{B}}_{+}\left(\phi, d f, A_{0}\right)  \tag{4-12}\\
\left(\partial_{t}+i|\xi|\right)(\hat{\phi}-i|\xi| \hat{f}) & =-\hat{\mathscr{B}}_{-}\left(\phi, d f, A_{0}\right) \tag{4-13}
\end{align*}
$$

where

$$
\begin{equation*}
\hat{\mathscr{B}}_{ \pm}=-\left[\widehat{\partial_{1} f, \partial_{2}} f\right]+\left[\widehat{A_{0}, \phi}\right] \pm\left(\frac{\xi_{1}}{|\xi|}\left[\widehat{\partial_{2} f, \phi}\right]-\frac{\xi_{2}}{|\xi|}\left[\widehat{\partial_{1} f, \phi}\right]+\frac{\xi_{j}}{|\xi|}\left[\widehat{A_{0}, \partial_{j}} f\right]\right) . \tag{4-14}
\end{equation*}
$$

Indeed, multiply (4-10) by $i /|\xi|$, and first add the resulting equation to (4-11) to obtain (4-12), and then subtract it from (4-11) to obtain (4-13). Now we let

$$
\begin{equation*}
\hat{\phi}+i|\xi| \hat{f}=\left(\partial_{t}+i|\xi|\right) \hat{u} \quad \text { and } \quad \hat{\phi}-i|\xi| \hat{f}=\left(\partial_{t}-i|\xi|\right) \hat{v}, \tag{4-15}
\end{equation*}
$$

where $u, v: \mathbb{R}^{2+1} \rightarrow \mathfrak{g}$. This gives

$$
\square u=\mathscr{B}_{+}\left(\phi, d f, A_{0}\right), \quad \square v=\mathscr{B}_{-}\left(\phi, d f, A_{0}\right) .
$$

See Remark 4.1 below.
Now we discuss initial data. From (4-15), we have

$$
\begin{gather*}
\partial_{t} \widehat{u(0)}=\hat{\phi}_{0}+i|\xi| \widehat{f(0)}-i|\xi| \widehat{u_{(0)}},  \tag{4-16}\\
\partial_{t} \widehat{v(0)}=\hat{\phi}_{0}-i|\xi| \widehat{f(0)}+i|\xi| \widehat{v(0)} . \tag{4-17}
\end{gather*}
$$

We are free to choose any data for $u$ and $v$, as long as in the end we can recover the original data for $\phi$ and $A$. Hence we just let $u(0)=v(0)=0$. We still need to say what $|\xi| \widehat{f(0)}$ is in terms of the initial data $\left(a, \phi_{0}\right)$. Let $\hat{h}=|\xi| \widehat{f(0)}$. By (4-1) and (4-5) we ahve $a_{1}=A_{1}(0)=-\partial_{2} f(0)$ and $a_{2}=A_{2}(0)=\partial_{1} f(0)$. Therefore

$$
R_{1} h=a_{2}, \quad R_{2} h=-a_{1},
$$

where $R_{j}$ denotes the Riesz transform, $(-\Delta)^{-1 / 2} \partial_{j}$. Differentiate the first equation with respect to $x$, the second with respect to $y$, and add them together to obtain

$$
\Delta D^{-1} h=\partial_{1} a_{2}-\partial_{2} a_{1}
$$

So

$$
\begin{equation*}
h=R_{2} a_{1}-R_{1} a_{2} \tag{4-18}
\end{equation*}
$$

It follows that the initial data for $u$ and $v$ are

$$
\begin{equation*}
u(0)=v(0)=0, \quad \partial_{t} u(0)=\phi_{0}+i h, \quad \partial_{t} v(0)=\phi_{0}-i h \tag{4-19}
\end{equation*}
$$

with $h$ defined by (4-18).
In summary, the monopole equation in the Coulomb gauge

$$
F_{A}=* D_{A} \phi, \quad d^{*} A=0,
$$

with initial data (4-1) can be rewritten as the system

$$
\begin{align*}
\square u & =\mathscr{B}_{+}\left(\phi, \nabla f, A_{0}\right), \\
\square v & =\mathscr{B}_{-}\left(\phi, \nabla f, A_{0}\right),  \tag{aME}\\
\triangle A_{0} & =\mathscr{C}_{( }\left(\phi, \nabla f, A_{0}\right),
\end{align*}
$$

where

$$
\begin{align*}
\mathscr{C} & =-\partial_{1}\left[A_{0}, \partial_{2} f\right]+\partial_{2}\left[A_{0}, \partial_{1} f\right]+\partial_{j}\left[\partial_{j} f, \phi\right],  \tag{4-20}\\
\mathscr{B}_{ \pm} & =-\mathscr{B}_{1} \mp i \mathscr{B}_{2}+\mathscr{B}_{3} \mp i \mathscr{B}_{4}, \tag{4-21}
\end{align*}
$$

and

$$
\begin{equation*}
\mathscr{B}_{1}=\left[\partial_{1} f, \partial_{2} f\right], \quad \mathscr{B}_{2}=R_{1}\left[\partial_{2} f, \phi\right]-R_{2}\left[\partial_{1} f, \phi\right], \quad \mathscr{B}_{3}=\left[A_{0}, \phi\right], \quad \mathscr{B}_{4}=R_{j}\left[A_{0}, \partial_{j} f\right] . \tag{4-22}
\end{equation*}
$$

The initial data for (aME) is given by (4-19).
Remark 4.1. $u$ and $v$ are our new unknowns, but we are really interested in $\phi$ and $d f$. Therefore, we observe that once we know what $u$ and $v$ are, we can determine $\phi$ and $d f$ by using

$$
\begin{equation*}
\hat{\phi}=\frac{\left(\partial_{t}+i|\xi|\right) \hat{u}+\left(\partial_{t}-i|\xi|\right) \hat{v}}{2}, \quad i|\xi| \hat{f}=\frac{\left(\partial_{t}+i|\xi|\right) \hat{u}-\left(\partial_{t}-i|\xi|\right) \hat{v}}{2} \tag{4-23}
\end{equation*}
$$

or equivalently

$$
\begin{equation*}
\phi=\frac{\left(\partial_{t}+i D\right) u+\left(\partial_{t}-i D\right) v}{2}, \quad \partial_{j} f=-i R_{j}\left(\frac{\left(\partial_{t}+i D\right) u-\left(\partial_{t}-i D\right) v}{2}\right) \tag{4-24}
\end{equation*}
$$

From $d f$ we get $A$ by letting $A=* d f$. Finally, for simplicity we usually keep the nonlinearities in terms of $\phi$ and $d f$. However, since $\phi$ and $d f$ can be written in terms of derivatives of $u$ and $v$ we sometimes write $\mathscr{B}_{ \pm}\left(\phi, d f, A_{0}\right)$ as $\mathscr{B}_{ \pm}\left(\partial u, \partial v, A_{0}\right)$.
Remark 4.2. (aME) has some resemblance to a system considered by Selberg [2002a] for the Maxwell-Klein-Gordon (MKG) equations, where he successfully obtains almost optimal local wellposedness in dimensions $1+4$. Besides the dimension considered, there are two fundamental technical differences applicable to our problem. First comes from the fact that the monopole equation we consider here is an example of a system in the nonabelian gauge theory whereas MKG is an example of a system in the abelian gauge theory. The existence of a global Coulomb gauge requires smallness of initial data in the nonabelian gauge theories, but is not needed in the abelian theories. Another technical difference arises from Selberg being able to solve the elliptic equation for his temporal variable $A_{0}$ using the Riesz representation theorem, where he does not require smallness of the initial data. The elliptic equation in (aME) is more difficult, so we include $A_{0}$ in the Picard iteration. As a result we are not able to allow large data by taking a small time interval, which we could do if we only had the two wave equations. Finally, we point out that the proof of our estimates involving $A_{0}$ is modeled after Selberg's proof in [Selberg 2002a] (see Remark 5.1 and Section 5C).

4B. Return to the monopole equation. Now we have a theorem, where we show how LWP for (aME) implies LWP for (ME) in the Coulomb gauge. For completeness, we first state exactly what we mean by LWP of (aME).

Let $r \in(0, \min (2 s, 1+s)], s>0$. Consider the system (aME) with initial data

$$
\left.\left(u, u_{t}\right)\right|_{t=0}=\left(u_{0}, u_{1}\right) \quad \text { and }\left.\quad\left(v, v_{t}\right)\right|_{t=0}=\left(v_{0}, v_{1}\right)
$$

in $H^{s+1} \times H^{s}$, then (aME) is LWP if the following conditions are satisfied:
Local existence. There exist $T>0$ depending continuously on the norm of the initial data, and functions

$$
A_{0} \in C_{b}\left([0, T], \dot{H}^{r}\right), \quad u, v \in \mathscr{H}_{T}^{s+1, \theta} \hookrightarrow C_{b}\left([0, T], H^{s+1}\right) \cap C_{b}^{1}\left([0, T], H^{s}\right)
$$

which solve (aME) on $[0, T] \times \mathbb{R}^{2}$ in the sense of distributions and such that the initial conditions are satisfied.
Uniqueness. If $T>0$ and $\left(A_{0}, u, v\right)$ and $\left(A_{0}^{\prime}, u^{\prime}, v^{\prime}\right)$ are two solutions of $(\mathrm{aME})$ on $(0, T) \times \mathbb{R}^{2}$ belonging to

$$
C_{b}\left([0, T], \dot{H}^{r}\right) \times \mathscr{H}_{T}^{s+1, \theta} \times \mathscr{H}_{T}^{s+1, \theta}
$$

with the same initial data, then $\left(A_{0}, u, v\right)=\left(A_{0}^{\prime}, u^{\prime}, v^{\prime}\right)$ on $(0, T) \times \mathbb{R}^{2}$.
Continuous dependence on initial data. For any $\left(u_{0}, u_{1}\right),\left(v_{0}, v_{1}\right) \in H^{s+1} \times H^{s}$ there is a neighborhood $U$ of the initial data such that the solution map $\left(u_{0}, u_{1}\right),\left(v_{0}, v_{1}\right) \rightarrow\left(A_{0}, u, v\right)$ is continuous from $U$ into $C_{b}\left([0, T], \dot{H}^{r}\right) \times\left(C_{b}\left([0, T], H^{s+1}\right) \cap C_{b}^{1}\left([0, T], H^{s}\right)\right)^{2}$.

In fact, by the results in [Selberg 2002b] combined with estimates for the elliptic equation, we can show the stronger estimates

$$
\begin{align*}
\left\|u-u^{\prime}\right\|_{\mathscr{K}_{T}^{s+1, \theta}}+\left\|v-v^{\prime}\right\|_{\mathscr{H}_{T}^{s+1, \theta}} & +\left\|A_{0}-A_{0}^{\prime}\right\|_{C_{b}\left([0, T], \dot{H}^{r}\right)} \\
& \lesssim\left\|u_{0}-u_{0}^{\prime}\right\|_{H^{s+1}}+\left\|u_{1}-u_{1}^{\prime}\right\|_{H^{s}}+\left\|v_{0}-v_{0}^{\prime}\right\|_{H^{s+1}}+\left\|v_{1}-v_{1}^{\prime}\right\|_{H^{s}} \tag{4-25}
\end{align*}
$$

where $\left(u_{0}^{\prime}, u_{1}^{\prime}\right),\left(v_{0}^{\prime}, v_{1}^{\prime}\right)$ are sufficiently close to $\left(u_{0}, u_{1}\right),\left(v_{0}, v_{1}\right)$.
Remark 4.3. Note that below we have no restriction on $s$, that is, if we could show (aME) is LWP in $H^{s+1} \times H^{s}, s>0$, we would get LWP of (ME) in the Coulomb gauge in $H^{s}$ for $s>0$ as well.

Theorem 4.1. Consider (ME) in the Coulomb gauge with the following initial data in $H^{s}$ for $s>0$ :

$$
\begin{equation*}
\left.A_{i}\right|_{t=0}=a_{i}, \quad i=1,2,\left.\quad \phi\right|_{t=0}=\phi_{0}, \quad \text { with } \partial^{i} a_{i}=0 \tag{4-26}
\end{equation*}
$$

Then local wellposedness of (aME) with initial data as in (4-19) implies local wellposedness of (ME) in the Coulomb gauge with initial data given by (4-26).
Proof. First, in view of Section 4A it is clear that if $u, v$ satisfy (aME) with initial data as in (4-19), then solutions of (ME) in the Coulomb gauge satisfy the initial data as given in (4-26).
Local existence. From (4-24), if

$$
u, v \in \mathscr{H}_{T}^{s+1, \theta}, \quad \text { then } \quad \phi, A=* d f \in H_{T}^{s, \theta},
$$

as needed. We now verify that if ( $u, v, A_{0}$ ) solve (aME), then ( $\phi, d f, A_{0}$ ) solve (ME) in the Coulomb gauge. Since $A=* d f, A$ is in the Coulomb gauge: $d * A=-* d *(* d f)=0$. Next note that (ME) in the

Coulomb gauge is equivalent to (4-7) and (4-8). Suppose $u, v, A_{0}$ solve (aME). It follows ( $d f, \phi$ ) solve (4-12) and (4-13). Add (4-12) to (4-13) to recover (4-11), which is equivalent to (4-8).

Next given (aME) we need to show (4-7) holds. Write (4-7) in coordinates,

$$
\begin{align*}
\partial_{1} A_{0}-\partial_{2} \phi+\partial_{t} \partial_{2} f & =\left[\partial_{1} f, \phi\right]-\left[A_{0}, \partial_{2} f\right],  \tag{4-27}\\
\partial_{2} A_{0}+\partial_{1} \phi-\partial_{t} \partial_{1} f & =\left[\partial_{2} f, \phi\right]+\left[A_{0}, \partial_{1} f\right] . \tag{4-28}
\end{align*}
$$

From the elliptic equation in (aME) we have:

$$
\begin{equation*}
A_{0}=\Delta^{-1}\left(-\partial_{1}\left[A_{0}, \partial_{2} f\right]+\partial_{2}\left[A_{0}, \partial_{1} f\right]+\partial_{1}\left[\partial_{1} f, \phi\right]+\partial_{2}\left[\partial_{2} f, \phi\right]\right) \tag{4-29}
\end{equation*}
$$

Also subtract (4-12) from (4-13) and multiply by $|\xi|$ on both sides to obtain (4-9), which implies

$$
\begin{equation*}
\phi-\partial_{t} f=\Delta^{-1}\left(\partial_{j}\left[A_{0}, \partial_{j} f\right]-\partial_{2}\left[\partial_{1} f, \phi\right]+\partial_{1}\left[\partial_{2} f, \phi\right]\right) \tag{4-30}
\end{equation*}
$$

In order to recover (4-27), first use (4-29) to get

$$
\begin{equation*}
\partial_{1} A_{0}=\Delta^{-1}\left(-\partial_{1}^{2}\left[A_{0}, \partial_{2} f\right]+\partial_{1} \partial_{2}\left[A_{0}, \partial_{1} f\right]+\partial_{1}^{2}\left[\partial_{1} f, \phi\right]+\partial_{1} \partial_{2}\left[\partial_{2} f, \phi\right]\right) \tag{4-31}
\end{equation*}
$$

Next use (4-30) to get

$$
\begin{equation*}
\partial_{2}\left(\phi-\partial_{t} f\right)=\Delta^{-1}\left(\partial_{2} \partial_{1}\left[A_{0}, \partial_{1} f\right]+\partial_{2}^{2}\left[A_{0}, \partial_{2} f\right]-\partial_{2}^{2}\left[\partial_{1} f, \phi\right]+\partial_{2} \partial_{1}\left[\partial_{2} f, \phi\right]\right) \tag{4-32}
\end{equation*}
$$

and subtract this from (4-31) to get (4-27) as needed. We recover (4-28) in the exactly same way.
Continuous dependence on initial data. We would like to show that

$$
\begin{align*}
\left\|A_{0}-A_{0}^{\prime}\right\|_{C_{b}\left([0, T], \dot{H}^{r}\right)}+\left\|A_{1}-A_{1}^{\prime}\right\|_{H_{T}^{s, \theta}}+\| A_{2}- & A_{2}^{\prime}\left\|_{H_{T}^{s, \theta}}+\right\| \phi-\phi^{\prime} \|_{H_{T}^{s, \theta}} \\
& \lesssim\left\|a_{1}-a_{1}^{\prime}\right\|_{H^{s}}+\left\|a_{2}-a_{2}^{\prime}\right\|_{H^{s}}+\left\|\phi_{0}-\phi_{0}^{\prime}\right\|_{H^{s}} \tag{4-33}
\end{align*}
$$

for any $a_{1}^{\prime}, a_{2}^{\prime}, \phi_{0}^{\prime}$ sufficiently close to $a_{1}, a_{2}, \phi_{0}$. In view of LWP for (aME) with data given by

$$
u(0)=v(0)=0, \quad \partial_{t} u(0)=\phi_{0}+i h, \quad \partial_{t} v(0)=\phi_{0}-i h, \quad h=R_{2} a_{1}-R_{1} a_{2}
$$

and by (4-25) we have

$$
\begin{align*}
\left\|u-u^{\prime}\right\|_{\mathscr{H}_{T}^{s+1, \theta}}+\left\|v-v^{\prime}\right\|_{\mathscr{H}_{T}^{s+1, \theta}} & +\left\|A_{0}-A_{0}^{\prime}\right\|_{C_{b}\left([0, T], \dot{H}^{r}\right)} \\
& \lesssim\left\|u_{0}^{\prime}\right\|_{H^{s+1}}+\left\|\phi_{0}+i h-u_{1}^{\prime}\right\|_{H^{s}}+\left\|v_{0}^{\prime}\right\|_{H^{s+1}}+\left\|\phi_{0}-i h-v_{1}^{\prime}\right\|_{H^{s}} \tag{4-34}
\end{align*}
$$

for all $u_{0}^{\prime}, v_{0}^{\prime}, u_{1}^{\prime}, v_{1}^{\prime}$ satisfying

$$
\begin{equation*}
\left\|u_{0}^{\prime}\right\|_{H^{s+1}}+\left\|\phi_{0}+i h-u_{1}^{\prime}\right\|_{H^{s}}+\left\|v_{0}^{\prime}\right\|_{H^{s+1}}+\left\|\phi_{0}-i h-v_{1}^{\prime}\right\|_{H^{s}} \leq \delta, \quad \text { for some } \delta>0 \tag{4-35}
\end{equation*}
$$

In particular choose

$$
\begin{equation*}
u_{0}^{\prime}=v_{0}^{\prime}=0, \quad u_{1}^{\prime}=\phi_{0}^{\prime}+i h^{\prime}, \quad v_{1}^{\prime}=\phi_{0}^{\prime}-i h^{\prime}, \quad h^{\prime}=R_{2} a_{1}^{\prime}-R_{1} a_{2}^{\prime} \tag{4-36}
\end{equation*}
$$

such that

$$
\begin{align*}
\left\|\phi_{0}+i h-\phi_{0}^{\prime}-i h^{\prime}\right\|_{H^{s}}+\left\|\phi_{0}-i h-\phi_{0}^{\prime}+i h^{\prime}\right\|_{H^{s}} & \lesssim\left\|\phi_{0}-\phi_{0}^{\prime}\right\|_{H^{s}}+\left\|R_{1}\left(a_{2}-a_{2}^{\prime}\right)\right\|_{H^{s}}+\left\|R_{2}\left(a_{1}-a_{1}^{\prime}\right)\right\|_{H^{s}}  \tag{4-37}\\
& \leq\left\|\phi_{0}-\phi_{0}^{\prime}\right\|_{H^{s}}+\left\|a_{1}-a_{1}^{\prime}\right\|_{H^{s}}+\left\|a_{2}-a_{2}^{\prime}\right\|_{H^{s}} \\
& \leq \delta .
\end{align*}
$$

Then, by (4-34)-(4-37), $\left\|A_{0}-A_{0}^{\prime}\right\|_{C_{b}\left([0, T], \dot{H}^{r}\right)}$ is bounded by the right side of (4-33). Next observe that

$$
\begin{aligned}
\left\|A_{1}-A_{1}^{\prime}\right\|_{H_{T}^{s, \theta}} & \lesssim\left\|R_{2}\left(\partial_{t}+i D\right)\left(u-u^{\prime}\right)\right\|_{H_{T}^{s, \theta}}+\left\|R_{2}\left(\partial_{t}-i D\right)\left(v-v^{\prime}\right)\right\|_{H_{T}^{s, \theta}} \\
& \leq\left\|u-u^{\prime}\right\|_{\mathscr{H}_{T}^{s+1, \theta}}+\left\|v-v^{\prime}\right\|_{\mathscr{H}_{T}^{s+1, \theta}} .
\end{aligned}
$$

So again by (4-34)-(4-37) \| $A_{1}-A_{1}^{\prime} \|_{H_{T}^{s, \theta}}$ is bounded by the right side of (4-33). We bound the difference for $A_{2}$ and $\phi$ in a similar fashion.
Uniqueness. By LWP of (aME), $A_{0}$ is unique in the required class. We need to show $A$ and $\phi$ are unique in $H_{T}^{s, \theta}$. However, by (4-33) this is obvious.

## 5. Proof of the Main Theorem

By Theorem 4.1 it is enough to show LWP for (aME). We start by explaining how we are going to perform our iteration.

5A. Set up of the iteration. Equations (aME) are written for functions $u$ and $v$. Nevertheless, functions $u$ and $v$ are only our auxiliary functions, and we are really interested in solving for $d f$ and $\phi$. In addition, the nonlinearities $\mathscr{B}_{ \pm}$are a linear combination of terms $\mathscr{B}_{i}, i=1,2,3,4$, given by (4-22), and the $\mathscr{B}_{i}$ are written in terms of $\phi, d f$ and $A_{0}$. Also, when we do our estimates, it is easier to keep the $\mathscr{B}_{i}$ in terms of $\phi$ and $d f$ with the exception of $\mathscr{B}_{2}$, which we rewrite in terms of $\partial u$ and $\partial v$ (see Section 5B2 for the details). These comments motivate the following procedure for our iteration. Start with $\phi_{-1}=d f_{-1}=0$. Then $\mathscr{B}_{ \pm} \equiv 0$. Solve the homogeneous wave equations for $u_{0}, v_{0}$ with the initial data given by (4-19). Then to solve for $d f_{0}, \phi_{0}$, use (4-24). Then feed $\phi_{0}$ and $d f_{0}$ into the elliptic equation,

$$
\begin{equation*}
\Delta A_{0,0}=-d^{*}\left(\left[A_{0,0}, * d f_{0}\right]+\left[d f_{0}, \phi_{0}\right]\right) \tag{5-1}
\end{equation*}
$$

and solve for $A_{0,0}$. Next we take $d f_{0}, \phi_{0}$ and $A_{0,0}$ plug them into $\mathscr{B}_{1}, \mathscr{B}_{3}, \mathscr{B}_{4}$, but rewrite $\mathscr{B}_{2}$ in terms of $\partial u_{0}, \partial v_{0}$. We continue in this manner, so at the $j$ th step of the iteration, $j \geq 1$, we solve

$$
\begin{aligned}
\square u_{j} & =-\mathscr{B}_{1}\left(\nabla f_{j-1}\right)-i \mathscr{B}_{2}\left(\partial u_{j-1}, \partial v_{j-1}\right)+\mathscr{B}_{3}\left(A_{0, j-1}, \phi_{j-1}\right)-i \mathscr{B}_{4}\left(A_{0, j-1}, \nabla f_{j-1}\right), \\
\square v_{j} & =-\mathscr{B}_{1}\left(\nabla f_{j-1}\right)+i \mathscr{B}_{2}\left(\partial u_{j-1}, \partial v_{j-1}\right)+\mathscr{B}_{3}\left(A_{0, j-1}, \phi_{j-1}\right)+i \mathscr{B}_{4}\left(A_{0, j-1}, \nabla f_{j-1}\right), \\
\triangle A_{0, j} & =-d^{*}\left(\left[A_{0, j}, * d f_{j}\right]+\left[d f_{j}, \phi_{j}\right]\right) .
\end{aligned}
$$

5B. Estimates needed. The elliptic equation is discussed in Section 5C. Therefore we begin by discussing the inversion of the wave operator in $\mathscr{H}^{s+1, \theta}$ spaces. The main idea is that for the purposes of local in time estimates $\square^{-1}$ can be replaced with $\Lambda_{+}^{-1} \Lambda_{-}^{-1}$. The first estimates, leading to wellposedness for small initial data, were proved by Klainerman and Machedon [1995]. The small data assumption was removed by Selberg [2002b], who showed that by introducing $\varepsilon$ small enough in the invertible version
of the wave operator, that is, $\Lambda_{+}^{-1} \Lambda_{-}^{-1+\varepsilon}$, we can use initial data as large as we wish. ${ }^{5}$ Selberg [2002b] also gave a very useful, general framework for local wellposedness of wave equations, which reduces the proof of the Main Theorem to establishing the estimates below, for the nonlinearities $\mathscr{B}_{ \pm}$, and to combining them with appropriate elliptic estimates from Section 5C. The needed estimates for $\mathscr{B}_{ \pm}$are

$$
\begin{gather*}
\left\|\Lambda_{+}^{-1} \Lambda_{-}^{-1+\varepsilon} \mathscr{B}_{ \pm}\left(\partial u, \partial v, A_{0}\right)\right\|_{\mathscr{H}^{s+1, \theta}}  \tag{5-2}\\
\lesssim\|u\|_{\mathscr{H}_{s+1, \theta}}+\|v\|_{\mathscr{H}^{s+1, \theta}},  \tag{5-3}\\
\left\|\Lambda_{+}^{-1} \Lambda_{-}^{-1+\varepsilon}\left(\mathscr{B}_{ \pm}\left(\partial u, \partial v, A_{0}\right)-\mathscr{B}_{ \pm}\left(\partial u^{\prime}, \partial v^{\prime}, A_{0}^{\prime}\right)\right)\right\|_{\mathscr{H}_{s}+1, \theta} \\
\lesssim u-u^{\prime}\left\|_{\mathscr{H}^{s+1, \theta}}+\right\| v-v^{\prime} \|_{\mathscr{H}^{s+1, \theta}},
\end{gather*}
$$

where the suppressed constants depend continuously on the $\mathscr{H}^{s+1, \theta}$ norms of $u, u^{\prime}, v, v^{\prime}$. Since $\mathscr{B}_{ \pm}$are bilinear, (5-3) can follow from (5-2). In this paper small initial data is necessary (see Theorem 3.3 and Section 5C), so we do not need $\varepsilon$, but we keep it to make the estimates general. Let $\frac{1}{4}<s<\frac{1}{2}$ and set $\theta, \varepsilon$ as follows:

$$
\frac{3}{4}-\frac{\varepsilon}{2}<\theta \leq s+\frac{1}{2}-\varepsilon, \quad \text { and } \quad \theta<1-\varepsilon, \quad 0 \leq \varepsilon<\min \left(2 s-\frac{1}{2}, \frac{1}{2}\right)
$$

Next observe $\Lambda_{+} \Lambda_{-}^{1-\varepsilon} \mathscr{H}^{s+1, \theta}=H^{s, \theta-1+\varepsilon}$, as well as that

$$
\|\nabla f\|_{H^{s, \theta}},\|\phi\|_{H^{s, \theta}} \lesssim\|u\|_{\mathscr{H}_{s}+1, \theta}+\|v\|_{\mathscr{H} s+1, \theta} .
$$

Therefore, using (4-21) and (4-22), it is enough to prove the following:

$$
\begin{align*}
& \left\|\mathscr{B}_{1}\right\|_{H^{s, \theta-1+\varepsilon}}=\left\|\left[\partial_{1} f, \partial_{2} f\right]\right\|_{H^{s, \theta-1+\varepsilon}} \lesssim\|\nabla f\|_{H^{s, \theta}}^{2},  \tag{5-4}\\
& \left\|\mathscr{B}_{2}\right\|_{H^{s, \theta-1+\varepsilon}} \lesssim\left\|\left[\partial_{j} f, \phi\right]\right\|_{H^{s, \theta-1+\varepsilon}} \lesssim\left\|\partial_{j} f\right\|_{H^{s, \theta}}\|\phi\|_{H^{s, \theta}}  \tag{5-5}\\
& \left\|\mathscr{B}_{3}\right\|_{H^{s, \theta-1+\varepsilon}} \lesssim\left\|A_{0} \phi\right\|_{H^{s, \theta-1+\varepsilon}} \lesssim\left\|A_{0}\right\|\|\phi\|_{H^{s, \theta}},  \tag{5-6}\\
& \left\|\mathscr{B}_{4}\right\|_{H^{s, \theta-1+\varepsilon}} \lesssim\left\|A_{0} \partial_{j} f\right\|_{H^{s, \theta-1+\varepsilon}} \lesssim\left\|A_{0}\right\|\left\|\partial_{j} f\right\|_{H^{s, \theta}} \quad \text { for } j=1,2, \tag{5-7}
\end{align*}
$$

where the norm used for $A_{0}$ is immaterial, mainly because in Section 5C we show that

$$
\begin{equation*}
\left\|A_{0}\right\| \lesssim\|\nabla f\|_{H^{s, \theta}}\|\phi\|_{H^{s, \theta}} . \tag{5-8}
\end{equation*}
$$

A few remarks are in order. Estimate (5-4) corresponds to estimates for the null form $Q_{i j}$, and estimate (5-5) gives rise to a new null form $Q$ (this is discussed in the next two sections). $A_{0}$ in estimates (5-6) and (5-7) solves the elliptic equation in (aME), which results in a quite good regularity for $A_{0}$. As a result, we do not have to look for any special structures to get (5-6) and (5-7) to hold, so we can drop the brackets, and also treat these estimates as equivalent since $\phi$ and $d f$ exhibit the same regularity. Finally, since Riesz transforms are clearly bounded on $L^{2}$, we ignore them in the estimates needed in (5-5) and (5-7). The estimates (5-4) and (5-5) for the null forms are the most interesting. Hence we discuss them first, and then we consider the elliptic terms.

5B1. Null forms: proof of estimate (5-4). $\left[\partial_{1} f, \partial_{2} f\right]$ has a structure of a null form $Q_{i j}$ :

$$
\left[\partial_{1} f, \partial_{2} f\right]=\partial_{1} f \partial_{2} f-\partial_{2} f \partial_{1} f=Q_{12}(f, f)
$$

[^3]It follows that (5-4) is equivalent to

$$
\left\|Q_{12}(f, f)\right\|_{H^{s, \theta-1+\varepsilon}} \lesssim\|\nabla f\|_{H^{s, \theta}}^{2} .
$$

Fortunately the hard work for null forms of type $Q_{\alpha, \beta}$ in two dimensions is already carried out by Zhou [1997]. His proof is done using spaces $N^{s+1, \theta}$ with the norm given by ${ }^{6}$

$$
\begin{equation*}
\|u\|_{N^{s+1, \theta}}=\left\|\Lambda_{+}^{s+1} \Lambda_{-}^{\theta} u\right\|_{L^{2}} . \tag{5-9}
\end{equation*}
$$

In his work $\theta=s+\frac{1}{2}$. We state Zhou's result.
Theorem [Zhou 1997]. Consider in $\mathbb{R}^{2+1}$ the space-time norms (5-9) and functions $\varphi, \psi$ defined on $\mathbb{R}^{2+1}$. The estimates

$$
\left\|Q_{\alpha \beta}(\varphi, \psi)\right\|_{N^{s, s-1 / 2}} \lesssim\|\varphi\|_{N^{s+1, s+1 / 2}}\|\psi\|_{N^{s+1, s+1 / 2}}
$$

hold for any $\frac{1}{4}<s<\frac{1}{2}$.
Our iteration is done using spaces $\mathscr{H}^{s+1, \theta}$. Inspection of Zhou's proof shows that it could be easily modified to be placed in the context of $\mathscr{H}^{s+1, \theta}$ spaces. However, even though our auxiliary functions' iterates $u_{j}$ and $v_{j}$ belong to $\mathscr{H}^{s+1, \theta}$, from (4-24) we only have

$$
\begin{equation*}
d f \in H^{s, \theta} \Rightarrow\left\|\Lambda^{s} \Lambda_{-}^{\theta} D f\right\|_{L^{2}\left(\mathbb{R}^{2+1}\right)}<\infty \tag{5-10}
\end{equation*}
$$

but again inspection of Zhou's proof shows we can still handle $Q_{12}(f, f)$ given only that (5-10) holds. Zhou's proof works for $\frac{1}{4}<s<\frac{1}{2}$. However, it motivates an alternate proof that uses $\mathscr{H}^{s+1, \theta}$ and works for all values of $s>\frac{1}{4}$. The proof is closely related to the original proof in [Zhou 1997], but on the surface it seems more concise. The reason for this is that we use Theorem F from [Klainerman and Selberg 2002], which involves all the technicalities. See [Czubak 2008] for the details.

5B2. Null forms: proof of estimate (5-5). We need

$$
\left\|\left[\partial_{j} f, \phi\right]\right\|_{H^{s, \theta-1+\varepsilon}} \lesssim\left\|\partial_{j} f\right\|_{H^{s, \theta}}\|\phi\|_{H^{s, \theta}}, \quad j=1,2
$$

However, analysis of the first iterate shows that for this estimate to hold we need $s>\frac{3}{4}$, so we need to work a little bit harder, and use (4-24) ${ }^{7}$

$$
\begin{equation*}
\left[\partial_{j} f, \phi\right]=-\frac{i}{4}\left[R_{j}\left(\partial_{t} u+i D u-\partial_{t} v+i D v\right), \partial_{t} u+i D u+\partial_{t} v-i D v\right] \tag{5-11}
\end{equation*}
$$

If we use the bilinearity of the bracket, we can group (5-11) by terms involving brackets of $u$ with itself, $v$ with itself, and then also by the terms that are mixed, that is, involve both $u$ and $v$. So we have

$$
\begin{aligned}
4 i\left[\partial_{j} f, \phi\right]= & {\left[R_{j}\left(\partial_{t}+i D\right) u,\left(\partial_{t}+i D\right) u\right]-\left[R_{j}\left(\partial_{t}-i D\right) v,\left(\partial_{t}-i D\right) v\right] } \\
& +\left[R_{j}\left(\partial_{t}+i D\right) u,\left(\partial_{t}-i D\right) v\right]-\left[R_{j}\left(\partial_{t}-i D\right) v,\left(\partial_{t}+i D\right) u\right] .
\end{aligned}
$$

Since $u$ and $v$ are matrix-valued and do not commute we need to combine the last two brackets to take advantage of a null form structure. This corresponds to (5-13) below (note the plus sign in the formula).

The needed estimates are contained in the following theorem.

[^4]Theorem 5.1. Suppose $s>\frac{1}{4}$ and $\frac{3}{4}-\frac{\varepsilon}{2}<\theta \leq s+\frac{1}{2}$, with $\theta<1-\varepsilon$ and $0 \leq \varepsilon<\min \left(2 s-\frac{1}{2}, \frac{1}{2}\right)$. Let $Q(\varphi, \psi)$ be given by

$$
\begin{align*}
& Q(\varphi, \psi)=\left(\partial_{t} \pm i D\right) R_{j} \varphi\left(\partial_{t} \pm i D\right) \psi-\left(\partial_{t} \pm i D\right) \varphi\left(\partial_{t} \pm i D\right) R_{j} \psi, \quad \text { or }  \tag{5-12}\\
& Q(\varphi, \psi)=\left(\partial_{t} \pm i D\right) R_{j} \varphi\left(\partial_{t} \mp i D\right) \psi+\left(\partial_{t} \pm i D\right) \varphi\left(\partial_{t} \mp i D\right) R_{j} \psi . \tag{5-13}
\end{align*}
$$

Then

$$
\begin{equation*}
Q\left(\mathscr{H}^{s+1, \theta}, \mathscr{H}^{s+1, \theta}\right) \hookrightarrow H^{s, \theta-1+\varepsilon} \tag{5-14}
\end{equation*}
$$

or, equivalently,

$$
\begin{equation*}
\|Q(\varphi, \psi)\|_{H^{s, \theta-1+\varepsilon}} \lesssim\|\varphi\|_{\mathscr{H}^{s}+1, \theta}\|\psi\|_{\mathscr{H}^{s+1, \theta}} \tag{5-15}
\end{equation*}
$$

Proof. We show the details only for

$$
\left(\partial_{t}+i D\right) R_{j} \varphi\left(\partial_{t}-i D\right) \psi+\left(\partial_{t}+i D\right) \varphi\left(\partial_{t}-i D\right) R_{j} \psi
$$

as the rest follows similarly. Observe that the symbol of $Q$ is then

$$
q(\tau, \xi, \lambda, \eta)=\left(\frac{\xi_{j}}{|\xi|}+\frac{\eta_{j}}{|\eta|}\right)(\tau+|\xi|)(\lambda-|\eta|)
$$

Suppose $\tau \lambda \geq 0$, then

$$
q \leq 2|(\tau+|\xi|)(\lambda-|\eta|)| \leq \begin{cases}2| | \tau|+|\xi||| | \lambda|-|\eta|| & \text { if } \quad \tau, \lambda \geq 0 \\ 2| | \tau|-|\xi||| | \eta|+|\lambda|| & \text { if } \quad \tau, \lambda \leq 0 .\end{cases}
$$

It follows

$$
\begin{equation*}
\iint_{\tau \lambda \geqq 0}\left|\Lambda^{s} \Lambda_{-}^{\theta-1+\varepsilon} Q(\varphi, \psi)\right|^{2} d \tau d \xi \lesssim\left\|D_{+} \varphi D_{-} \psi\right\|_{H^{s, \theta-1+\varepsilon}}^{2}+\left\|D_{-} \varphi D_{+} \psi\right\|_{H^{s, \theta-1+\varepsilon}}^{2} \tag{5-16}
\end{equation*}
$$

and the estimate follows by Theorem 5.2 below.
Suppose $\tau \lambda<0$. If we break down the computations into the region

$$
\begin{equation*}
\{(\tau, \xi),(\lambda, \eta):|\tau| \geq 2|\xi| \text { or }|\lambda| \geq 2|\eta|\} \tag{5-17}
\end{equation*}
$$

and its complement, then in the region (5-17), we bound $q$ by

$$
q \leq 2(|\tau|+|\xi|)(|\lambda|+|\eta|)
$$

since there we do not need any special structure. (This is a simple exercise in this region; see [Czubak 2008, Appendix B].)

In the complementary region, we have

$$
q \leq 4|\xi||\eta|\left|\frac{\xi_{i}}{|\xi|}+\frac{\eta_{i}}{|\eta|}\right|
$$

which is the absolute value of the symbol of the null form $Q_{t j}$ in the first iterate. It has received a lot of attention, but we have not seen a reference where it was discussed in a context other than that of the initial data in $H^{s+1} \times H^{s}$. This may be because it has not come up as a nonlinearity before, and/or because it can be handled in the same way as the null form $Q_{i j}$. The details are in [Czubak 2008].

Now we prove an estimate needed to show (5-16) is bounded by the square of the right side of (5-15).
Theorem 5.2. Let $s>0$, $\max \left(\frac{1}{2}, 1-s\right)<\theta<1$, and $0 \leq \varepsilon \leq 1-\theta$. Then

$$
\left\|D_{+} \varphi D_{-} \psi\right\|_{H^{s, \theta-1+\varepsilon}} \lesssim\|\varphi\|_{\mathscr{H}^{s+1, \theta}}\|\psi\|_{\mathscr{H}_{s}^{s+1, \theta}} .
$$

Proof. We would like to show that $\left\|\Lambda^{s} \Lambda_{-}^{\theta-1+\varepsilon}\left(D_{+} \varphi D_{-} \psi\right)\right\|_{L^{2}\left(\mathbb{R}^{2+1}\right)} \lesssim\|\varphi\|_{\mathscr{H} s+1, \theta}\|\psi\|_{\mathscr{H}^{s+1, \theta}}$. This follows from showing that

$$
H^{s, \theta} \cdot \mathscr{H}^{s+1, \theta-1} \hookrightarrow H^{s, \theta-1+\varepsilon},
$$

which by the product rule ${ }^{8}$ for the operator $\Lambda^{s}$ in turn follows from

$$
H^{0, \theta} \cdot \mathscr{H}^{s+1, \theta-1} \hookrightarrow H^{0, \theta-1+\varepsilon}, \quad H^{s, \theta} \cdot \mathscr{H}^{1, \theta-1} \hookrightarrow H^{0, \theta-1+\varepsilon} .
$$

It is easy to check that $\mathscr{H}^{s+1, \theta-1} \hookrightarrow H^{s+1+\theta-1,0}$ and $\mathscr{H}^{1, \theta-1} \hookrightarrow H^{\theta, 0}$, so we just need to show

$$
H^{0, \theta} \cdot H^{s+\theta, 0} \hookrightarrow H^{0, \theta-1+\varepsilon}, \quad H^{s, \theta} \cdot H^{\theta, 0} \hookrightarrow H^{0, \theta-1+\varepsilon},
$$

which are weaker than

$$
H^{0, \theta} \cdot H^{s+\theta, 0} \hookrightarrow L^{2}, \quad H^{s, \theta} \cdot H^{\theta, 0} \hookrightarrow L^{2}
$$

but those follow from the Klainerman-Selberg estimate (2-6) as long as $s+\theta>1$, which holds by the conditions we impose on $s$ and $\theta$.
An alternate approach could be to follow the set up used by [Klainerman and Machedon 1995] and estimate the integral directly.
5B3. Elliptic piece: proof of estimate (5-6). Recall we wish to show

$$
\begin{equation*}
\left\|A_{0} w\right\|_{H^{s, \theta-1+\varepsilon}} \lesssim\left\|A_{0}\right\|\|w\|_{H^{s, \theta}} \tag{5-18}
\end{equation*}
$$

We need this estimate during our iteration, so we really mean $A_{0, j}$, but for simplicity we omit writing of the index $j$. Now we choose a norm for $A_{0}$ to be anything that makes (5-18) possible to establish. This results in

$$
\left\|A_{0}\right\|=\left\|A_{0}\right\|_{L_{t}^{\tilde{p}} L_{x}^{\infty}}+\left\|D^{s} A_{0}\right\|_{L_{t}^{p} L_{x}^{q}}
$$

where

$$
\begin{equation*}
\tilde{p} \in\left(1-2 s, \frac{1}{2}\right), \quad \frac{2}{p}=1-\frac{1}{q}, \quad \max \left(\frac{1}{3}(1-2 s), \frac{s}{2}\right)<\frac{1}{q}<\frac{2}{3} s \tag{5-19}
\end{equation*}
$$

For now we assume we can show $A_{0} \in L_{t}^{\tilde{p}} L_{x}^{\infty} \cap L_{t}^{p} \dot{W}_{x}^{s, q}$ and delay the proof to Section 5 C , where the reasons for our choices of $\tilde{p}, p, q$ should become clear. We start by using $\theta-1+\varepsilon<0$ :

$$
\begin{equation*}
\left\|A_{0} w\right\|_{H^{s, \theta-1+\varepsilon}} \leq\left\|\Lambda^{s}\left(A_{0} w\right)\right\|_{L^{2}\left(\mathbb{R}^{2+1}\right)} \lesssim\left\|A_{0} w\right\|_{L^{2}\left(\mathbb{R}^{2+1}\right)}+\left\|D^{s}\left(A_{0} w\right)\right\|_{L^{2}\left(\mathbb{R}^{2+1}\right)} \tag{5-20}
\end{equation*}
$$

For the first term, by Hölder's inequality,

$$
\begin{array}{rlrl}
\left\|A_{0} w\right\|_{L^{2}\left(\mathbb{R}^{2+1}\right)} & \leq\left\|A_{0}\right\|_{L_{t}^{\tilde{p}} L_{x}^{\infty}}\|w\|_{L_{t}^{\tilde{p}^{\prime}} L_{x}^{2}}, & \text { for } \frac{1}{\tilde{p}}+\frac{1}{\tilde{p}^{\prime}}=\frac{1}{2}, \quad p \text { as in (5-19) } \\
& \lesssim\left\|A_{0}\right\|\|w\|_{H^{0, \theta}}, & & \text { by }(2-4)  \tag{5-21}\\
& \leq\left\|A_{0}\right\|\|w\|_{H^{s, \theta}} . &
\end{array}
$$

[^5]We bound the second term in (5-20) by
where

$$
\left\|D^{s}\left(A_{0} w\right)\right\|_{L^{2}\left(\mathbb{R}^{2+1}\right)} \lesssim \underbrace{\left\|A_{0}\right\|_{L_{t}^{\tilde{p}} L_{x}^{\infty}}\left\|D^{s} w\right\|_{L_{t}^{p^{\prime}} L_{x}^{2}}}_{I}+\underbrace{\left\|D^{s} A_{0}\right\|_{L_{t}^{p} L_{x}^{q}\|w\|_{L_{t}^{p^{\prime}} L_{x}^{q^{\prime}}}}, ~}_{I I}
$$

$$
\frac{1}{p}+\frac{1}{p^{\prime}}=\frac{1}{2}=\frac{1}{q}+\frac{1}{q^{\prime}}
$$

and $p, q$ are as in(5-19) and $\tilde{p}$ as in (5-21). I is handled similarly to (5-21) as follows. Apply (2-4) with $u=D^{s} w$ to obtain

$$
\begin{equation*}
I \lesssim\left\|A_{0}\right\|\left\|D^{s} w\right\|_{H^{0, \theta}} \leq\left\|A_{0}\right\|\|w\|_{H^{s, \theta}} . \tag{5-22}
\end{equation*}
$$

We now consider II. By the choices of $p, q$, the Klainerman-Selberg estimate (2-5) applies (see the discussion in Section 5D for an explanation) and gives

$$
\begin{equation*}
I I \leq\left\|A_{0}\right\|\|w\|_{L_{t}^{p^{\prime}} L_{x}^{q^{\prime}}} \lesssim\left\|A_{0}\right\|\|w\|_{H^{1-2 / q^{\prime}-1 / p^{\prime}, \theta}} . \tag{5-23}
\end{equation*}
$$

From (5-19) we also have

$$
\begin{equation*}
I I \lesssim\left\|A_{0}\right\|\|w\|_{H^{1-2 / q^{\prime}-1 / p^{\prime}, \theta}} \lesssim\left\|A_{0}\right\|\|w\|_{H^{s, \theta}} . \tag{5-24}
\end{equation*}
$$

and (5-18) follows now from (5-21), (5-22) and (5-24).
Remark 5.1. The above proof illustrates other difficulties due to working in two dimensions. Initially, we wanted to follow the proof of estimate (38) in [Selberg 2002a], and just use the $\left\|\Lambda^{s} A_{0}\right\|_{L_{t}^{p} L_{x}^{q}}$ norm. Unfortunately in two dimensions, the condition $s q>2$ needed to show that $A_{0} \in L_{t}^{p} L_{x}^{\infty}$ is disjoint from the conditions needed to use Klainerman-Tataru estimate (2-3) and establish that $\Lambda^{s} A_{0} \in L_{t}^{p} L_{x}^{q}$ in the first place. This resulted in the $L_{t}^{\tilde{p}} L_{x}^{\infty} \cap L_{t}^{p} \dot{W}_{x}^{s, q}$ space above and also having to employ the KlainermanSelberg estimate (2-5), which was not needed for the proof of [Selberg 2002a, estimate (38)].

5C. Elliptic regularity: estimates for $\boldsymbol{A}_{\mathbf{0}}$. Here we present a variety of a priori estimates for the nondynamical variable $A_{0}$. At each point we could add the index $j$ to $A_{0}, d f$ and $\phi$. Therefore the presentation also applies to the iterates $A_{0, j}$. It is an exercise to show that the estimates we obtain here are enough to solve for $A_{0, j}$ at each step as well as to close the iteration for $A_{0}$. Let $A_{0}$ solve

$$
\Delta A_{0}=-d^{*}\left[A_{0}, * d f\right]-d^{*}[d f, \phi]=-\partial_{1}\left[A_{0}, \partial_{2} f\right]+\partial_{2}\left[A_{0}, \partial_{1} f\right]+\partial_{j}\left[\partial_{j} f, \phi\right]
$$

There is a wide range of estimates $A_{0}$ satisfies. Nevertheless, the two spatial dimensions limit our "range of motion." For example, it does not seem possible to place $A_{0}(t)$ in $L^{2}$. We state the general results and only show the cases we need to prove $A_{0} \in L_{t}^{\tilde{p}} L_{x}^{\infty} \cap L_{t}^{p} \dot{W}_{x}^{s, q}$ as required in the last section. The rest of the cases can be found in [Czubak 2008]. We add that the proofs of both of the following theorems were originally inspired by the proof of estimate (45) in [Selberg 2002a]. We start with the homogeneous estimates.
Theorem 5.3. Let $s>0$, and let $0 \leq a \leq s+1$ be given, and suppose $1 \leq p \leq \infty$ and $1<q<\infty$ satisfy

$$
\begin{gather*}
\max \left(\frac{1}{3}(1+2 a-4 s), \frac{1}{2}(1+a-4 s), \frac{1}{2} \min (a, 1)\right)<\frac{1}{q}<\frac{1+a}{2},  \tag{5-25}\\
1-\frac{2}{q}+a-2 s \leq \frac{1}{p} \leq \frac{1}{2}\left(1-\frac{1}{q}\right), \quad \frac{1}{p}<\left(1-\frac{2}{q}+a\right) . \tag{5-26}
\end{gather*}
$$

(i) If $0 \leq a \leq 1$ and the $H^{s, \theta}$ norm of $\nabla f$ is sufficiently small, then $A_{0} \in L_{t}^{p} \dot{W}_{x}^{a, q}$ and we have the estimate

$$
\begin{equation*}
\left\|A_{0}\right\|_{L_{t}^{p} \dot{W}_{x}^{a, q}} \lesssim\|\phi\|_{H^{s, \theta}}\|\nabla f\|_{H^{s, \theta}} \tag{5-27}
\end{equation*}
$$

(ii) If $1<a \leq s+1$ and $A_{0} \in L_{t}^{p} L_{x}^{(1 / q-1 / 2)^{-1}}$, then $A_{0} \in L_{t}^{p} \dot{W}_{x}^{a, q}$ and we have

$$
\begin{equation*}
\left\|A_{0}\right\|_{L_{t}^{p} \dot{W}_{x}^{a, q}} \lesssim\left(\left\|A_{0}\right\|_{L_{t}^{p} L_{x}^{(1 / q-1 / 2)^{-1}}}+\|\phi\|_{H^{s, \theta}}\right)\|\nabla f\|_{H^{s, \theta}} . \tag{5-28}
\end{equation*}
$$

Corollary 5.4. Let $s>0$, then $A_{0} \in C_{b}\left(I: \dot{H}_{x}^{a}\right)$, where

$$
0<a \leq \begin{cases}2 s & \text { if } 0<s \leq 1 \\ 1+s & \text { if } 1<s\end{cases}
$$

Proof of Corollary 5.4. Suppose $0<s<\frac{1}{2}$. Then use part (i) of the theorem with $q=2$ and $p=\infty$ to obtain $A_{0} \in L_{t}^{\infty} \dot{H}_{x}^{a}$ for $a \leq 2 s . A_{0}$ continuous as a function of time easily follows from a contraction argument in $C_{b}\left(I: \dot{H}_{x}^{a}\right)$ using $L_{t}^{\infty} \dot{H}_{x}^{a}$ estimates. The case $s \geq \frac{1}{2}$ is considered in [Czubak 2008].

So far we just need $s>0$ in order to make the estimates work. The requirement for $s>\frac{1}{4}$ does not come in until we start looking at the nonhomogeneous spaces, where also the range of $p$ and $q$ is smaller. However, we can distinguish two cases, $a q<2$ and $a q>2$.

Theorem 5.5. Let $s>0$, and suppose the $H^{s, \theta}$ norm of $\nabla f$ is sufficiently small.
(i) If aq $<2$ for $0<a<\min (2 s, 1)$ and if $p$ and $q$ satisfy

$$
\begin{align*}
\max \left(\frac{1}{2}+a-2 s, \frac{a}{2}\right) & <\frac{1}{q}<\frac{1}{2}  \tag{5-29}\\
1-\frac{2}{q}+a-2 s & \leq \frac{1}{p}<\frac{1}{2}-\frac{1}{q} \tag{5-30}
\end{align*}
$$

then $A_{0} \in L_{t}^{p} W_{x}^{a, q}$ and we have the estimate

$$
\begin{equation*}
\left\|A_{0}\right\|_{L_{t}^{p} W_{x}^{a, q}} \lesssim\|\phi\|_{H^{s, \theta}}\|\nabla f\|_{H^{s, \theta}} \tag{5-31}
\end{equation*}
$$

(ii) If aq $>2$, we need $s>\frac{1}{4}$ and $0<a<\min (4 s-1,1+s, 2 s)$. Suppose $p$ and $q$ also satisfy

$$
\begin{gather*}
\max \left(\frac{a-s}{2}, \frac{1}{2}+a-2 s\right)<\frac{1}{q}<\frac{1}{2} \min (a, 1)  \tag{5-32}\\
1-\frac{2}{q}+a-2 s \leq \frac{1}{p}<\frac{1}{2}-\frac{1}{q} \tag{5-33}
\end{gather*}
$$

then $A_{0} \in L_{t}^{p} W_{x}^{a, q}$ and we have the estimate

$$
\begin{equation*}
\left\|A_{0}\right\|_{L_{t}^{p} W_{x}^{a, q}} \lesssim\|\phi\|_{H^{s, \theta}}\|\nabla f\|_{H^{s, \theta}} \tag{5-34}
\end{equation*}
$$

Corollary 5.6. If $s>\frac{1}{4}$ and the $H^{s, \theta}$ norm of $\nabla f$ is sufficiently small, we have in particular $A_{0} \in L_{t}^{p} L_{x}^{\infty}$ for $p$ satisfying

$$
\begin{equation*}
1-2 s<\frac{1}{p}<\frac{1}{2} \tag{5-35}
\end{equation*}
$$

and we have the estimate

$$
\begin{equation*}
\left\|A_{0}\right\|_{L_{t}^{p} L_{x}^{\infty}} \lesssim\|\phi\|_{H^{s, \theta}}\|\nabla f\|_{H^{s, \theta}} \tag{5-36}
\end{equation*}
$$

Proof of Corollary 5.6. For each $p \in\left(1-2 s, \frac{1}{2}\right)$ we can find some $a$ and $q$, which satisfy the conditions of Theorem 5.5, part (ii). The corollary then follows from the Sobolev embedding: $W^{a, q}\left(\mathbb{R}^{2}\right) \hookrightarrow L^{\infty}\left(\mathbb{R}^{2}\right)$ for $a q>2$.
Remark 5.2. Estimate (5-34) is where $s>\frac{1}{4}$ is needed. Conditions on $\frac{1}{p}$ in (5-33) are needed so we can use (below) the Klainerman-Tataru estimate (2-3). In order to be able to choose such $\frac{1}{p}$, obviously $1-\frac{2}{q}+a-2 s$ must be strictly less than $\frac{1}{2}-\frac{1}{q}$. This forces $\frac{1}{q}$ to be strictly greater than $\frac{1}{2}+a-2 s$. We also need $a q>2$ to use the Sobolev embedding in Corollary 5.6, so if we want to be able to find $q$ between $\frac{1}{2}+a-2 s$ and $\frac{a}{2}$, a is forced to be strictly less than $4 s-1$. Therefore $s$ must be greater than $\frac{1}{4}$. See below for another instance of requiring $s>\frac{1}{4}$.

5D. Proof of estimates needed in 5B3. Recall we would like to show $A_{0} \in L_{t}^{\tilde{p}} L_{x}^{\infty} \cap L_{t}^{p} \dot{W}_{x}^{s, q}$. Therefore, we are interested in part (i) of Theorem 5.3 and part ii) in Theorem 5.5, so we can conclude Corollary 5.6. Moreover, we need a specific case of part (i) in Theorem 5.3, because we need $A_{0} \in L_{t}^{p} \dot{W}_{x}^{s, q}$, where $p, q$ in addition satisfy

$$
\begin{equation*}
1-\frac{2}{p} \leq \frac{1}{q}<\frac{1}{2} \quad \text { and } \quad \frac{2}{q}-\frac{1}{2}+\frac{1}{p} \leq s \tag{5-37}
\end{equation*}
$$

so we can use the embedding

$$
\begin{equation*}
H^{s, \theta} \hookrightarrow H^{1-(1-2 / q)-(1 / 2-1 / p), \theta} \hookrightarrow L_{t}^{(1 / 2-1 / p)^{-1}} L_{x}^{(1 / 2-1 / q)^{-1}} \tag{5-38}
\end{equation*}
$$

in (5-23) and (5-24). When we put (5-37) together with (5-25) and (5-26) with $a=s$, we obtain the second line of (5-19), namely

$$
\begin{equation*}
\frac{2}{p}=1-\frac{1}{q}, \quad \max \left(\frac{1}{3}(1-2 s), \frac{s}{2}\right)<\frac{1}{q}<\frac{2}{3} s . \tag{5-39}
\end{equation*}
$$

Remark 5.3. Observe that in order to be able to find such $q$ we must have $s>\frac{1}{4}$.
Consider

$$
\begin{align*}
\left\|A_{0}\right\|_{L_{t}^{p}} \dot{W}_{x}^{s, q} & =\left\|\Delta^{-1}\left(d^{*}\left[A_{0}, * d f\right]+d^{*}[d f, \phi]\right)\right\|_{L_{t}^{p} \dot{W}_{x}^{s, q}} \\
& \lesssim\left\|D^{-1}\left(A_{0} \nabla f\right)\right\|_{L_{t}^{p}} \dot{W}_{x}^{s, q}+\left\|D^{-1}(\nabla f \phi)\right\|_{L_{t}^{p} \dot{W}_{x}^{s, q}} \\
& \lesssim\left\|D^{s-1}\left(A_{0} \nabla f\right)\right\|_{L_{t}^{p} L_{x}^{q}}+\left\|D^{s-1}(\nabla f \phi)\right\|_{L_{t}^{p} L_{x}^{q}}  \tag{5-40}\\
& \lesssim\left\|A_{0} \nabla f\right\|_{L_{t}^{p} L_{x}^{r}}+\left\|D^{s-1}(\nabla f \phi)\right\|_{L_{t}^{p} L_{x}^{q}},
\end{align*}
$$

where we use the Sobolev embedding with $\frac{1}{q}=\frac{1}{r}-\frac{1-s}{2}$. The latter term is bounded by $\|\nabla f\|_{H^{s, \theta}}\|\phi\|_{H^{s, \theta}}$ using the Klainerman-Tataru estimate (2-3), whose application we discuss in the section below. For the former we use $\frac{1}{r}=\frac{1}{q}+\frac{1-s}{2}=\left(\frac{1}{q}-\frac{s}{2}\right)+\frac{1}{2}$ :

$$
\begin{equation*}
\left\|A_{0} \nabla f\right\|_{L_{t}^{p} L_{x}^{r}} \leq\left\|A_{0}\right\|_{L_{t}^{p} L_{x}^{(1 / q-s / 2)^{-1}}}\|\nabla f\|_{L_{t}^{\infty} L_{x}^{2}} \lesssim\left\|A_{0}\right\|_{L_{t}^{p} \dot{W}_{x}^{s, q}}\|\nabla f\|_{H^{s, \theta}} \tag{5-41}
\end{equation*}
$$

Then if the $H^{s, \theta}$ norm of $\nabla f$ is sufficiently small, we obtain

$$
\begin{equation*}
\left\|A_{0}\right\|_{L_{t}^{p} \dot{W}_{x}^{s, q}} \lesssim\|\nabla f\|_{H^{s, \theta}}\|\phi\|_{H^{s, \theta}} \tag{5-42}
\end{equation*}
$$

as needed.
For the nonhomogeneous estimate, since here $\frac{1}{4}<s<\frac{1}{2}$ the upper bound for $a$ is simply $4 s-1$. In addition, for our purposes right now it suffices to show the estimate for one particular $a$. Therefore we set $0<a<\min (s, 4 s-1)$ for $\frac{1}{4}<s<\frac{1}{2}$, and we let $p, q$ satisfy (5-32) and (5-33). We have

$$
\begin{align*}
\left\|A_{0}\right\|_{L_{t}^{p} W_{x}^{a, q}} & \lesssim\left\|D^{-1}\left(A_{0} \nabla f\right)\right\|_{L_{t}^{p} W_{x}^{a, q}}+\left\|D^{-1}(\nabla f \phi)\right\|_{L_{t}^{p} W_{x}^{a, q}}  \tag{5-43}\\
& \lesssim\left\|D^{-1}\left(A_{0} \nabla f\right)\right\|_{L_{t}^{p} L_{x}^{q}}+\left\|D^{-1}(\nabla f \phi)\right\|_{L_{t}^{p} L_{x}^{q}}+\left\|D^{a-1}\left(A_{0} \nabla f\right)\right\|_{L_{t}^{p} L_{x}^{q}}+\left\|D^{a-1}(\nabla f \phi)\right\|_{L_{t}^{p} L_{x}^{q}} .
\end{align*}
$$

The Klainerman-Tataru estimate (2-3) handles the second and the last term (see below). Consider the first term:

$$
\begin{align*}
\left\|D^{-1}\left(A_{0} \nabla f\right)\right\|_{L_{t}^{p} L_{x}^{q}} & \lesssim\left\|A_{0} \nabla f\right\|_{L_{t}^{p} L_{x}^{p}}, \quad \text { for } \frac{1}{q}=\frac{1}{r}-\frac{1}{2}, \\
& \leq\left\|A_{0}\right\|_{L_{t}^{p} L_{x}^{q}\|\nabla f\|_{L_{t}^{\infty} L_{x}^{2}}} \\
& \leq\left\|A_{0}\right\|_{L_{t}^{p} W_{x}^{a, q}}\|\nabla f\|_{H^{s, \theta}} . \tag{5-44}
\end{align*}
$$

For the third term we have

$$
\begin{array}{rlr}
\left\|D^{a-1}\left(A_{0} \nabla f\right)\right\|_{L_{t}^{p} L_{x}^{q}} & \lesssim\left\|A_{0} \nabla f\right\|_{L_{t}^{p} L_{x}^{r}}, & \text { for } \frac{1}{q}=\frac{1}{r}-\frac{1-a}{2} \\
& \lesssim\left\|A_{0}\right\|_{L_{t}^{p} L_{x}^{q}}\left\|D^{a} \nabla f\right\|_{L_{t}^{\infty} L_{x}^{2}}, & \text { for } \frac{1}{r}=\frac{1}{q}+\left(\frac{1}{2}-\frac{a}{2}\right) \\
& \lesssim\left\|A_{0}\right\|_{L_{t}^{p} W_{x}^{a, q}}\|\nabla f\|_{H^{s, \theta}}, &
\end{array}
$$

Then as before, this completes the proof if the $H^{s, \theta}$ norm of $\nabla f$ is sufficiently small.
Applying the Klainerman-Tataru theorem. We said that several of the above estimates follow from the Klainerman-Tataru estimate (2-3). We need to check that this is in fact the case. We begin by stating the theorem. We state it for two dimensions only, and as it is given in [Klainerman and Selberg 2002] (the original result holds for $n \geq 2$ ).

Theorem [Klainerman and Tataru 1999]. Let $1 \leq p \leq \infty, 1 \leq q<\infty$. Assume that

$$
\begin{align*}
& \frac{1}{p} \leq \frac{1}{2}\left(1-\frac{1}{q}\right),  \tag{5-45}\\
& 0<\sigma<2\left(1-\frac{1}{q}-\frac{1}{p}\right),  \tag{5-46}\\
& s_{1}, s_{2}<1-\frac{1}{q}-\frac{1}{2 p},  \tag{5-47}\\
& s_{1}+s_{2}+\sigma=2\left(1-\frac{1}{q}-\frac{1}{2 p}\right) . \tag{5-48}
\end{align*}
$$

Then

$$
\left\|D^{-\sigma}(u v)\right\|_{L_{t}^{p} L_{x}^{q}\left(\mathbb{R}^{2}\right)} \lesssim\|u\|_{H^{s_{1}, \theta}}\|v\|_{H^{s_{2}, \theta}},
$$

provided $\theta>\frac{1}{2}$.

The first time we use the theorem is in (5-40) for the term $\left\|D^{s-1}(\nabla f \phi)\right\|_{L_{t}^{p} L_{x}^{q}}$. Note that $\sigma=1-s$. Clearly $1 \leq p \leq \infty, 1 \leq q<\infty$. Next by (5-39) $\frac{2}{p}=1-\frac{1}{q}$, so (5-45) holds. Since $s<\frac{1}{2}, \sigma>0$, and we can see (5-46) holds when we substitute $\frac{1}{2}-\frac{1}{2 q}$ for $\frac{1}{p}$ in the right side and use $\frac{1}{q}<\frac{2}{3} s$. Next we let $s_{1}=s_{2}$ and with $\sigma=1-s>0$, (5-48) implies (5-47), so we only check (5-48). To that end we must be able to choose $s_{1}$ so that

$$
2 s_{1}=1-\frac{2}{q}-\frac{1}{p}+s \leq 2 s
$$

which is equivalent to our condition on $p$ and one of the lower bounds on $\frac{1}{q}$.
The next place we use the theorem is in (5-43) for $\left\|D^{-1}(\nabla f \phi)\right\|_{L_{t}^{p} L_{x}^{q}},\left\|D^{a-1}(\nabla f \phi)\right\|_{L_{t}^{p} L_{x}^{q}}$, where $p$ and $q$ are as in (5-32) and (5-33) with $0<a<\min (s, 4 s-1)<1$. Then for $\sigma=1$, by the right side of (5-33), (5-46) holds and implies (5-45). Note, since (5-46) is true with $\sigma=1$, it is true with $\sigma=1-a$. Next, for $\sigma=1$ (5-48) gives (5-47) and also for $\sigma=1-a$ as long as $0<a<1$. So again it is sufficient to see we can have $s_{1}$ defined by (5-48) such that $s_{1} \leq s$, but for $\sigma=1-a$ that follows from the left side of (5-33), and shows we can find it for $\sigma=1$ as well.

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[^0]:    MSC2000: 35L70, 70S15.
    Keywords: monopole, null form, Coulomb gauge, wellposedness.

[^1]:    ${ }^{1}$ See Section 4 for the precise formula for $\mathscr{B}_{ \pm}$and $\mathscr{C}$.
    ${ }^{2}$ The discussion of the first iterate can be also found in the appendix of [Klainerman and Selberg 2002], and it can be deduced from the estimates and counterexamples found within [Foschi and Klainerman 2000].

[^2]:    ${ }^{3} \mathrm{We}$ are also going to employ a combination of the standard $L_{t}^{p} W_{x}^{s, q}$ spaces for $A_{0}$. See Section 5C.
    ${ }^{4}$ These spaces, together with results in [Selberg 2002b], allowed Klainerman and Selberg to present a unified approach to local wellposedness for wave maps, Yang-Mills and Maxwell-Klein-Gordon types of equations in [Klainerman and Selberg 2002], and are now the natural choice for low regularity subcritical local wellposedness for wave equations. See also [Tao 2006].

[^3]:    ${ }^{5}$ See also [Klainerman and Selberg 2002, Section 5] for an excellent discussion and motivation of the issues involved in the Picard iteration.

[^4]:    ${ }^{6}$ See [Selberg 1999, Section 3.5] for a comparison with $\mathscr{H}^{s+1, \theta}$ spaces.
    ${ }^{7}$ The obvious way is to just substitute for $\phi$ and leave $\partial_{j} f$ the same, but it is an exercise to see that this does not work - for several reasons!

[^5]:    ${ }^{8}$ On $L^{2}$ this is very easy to establish using triangle inequality. See [Klainerman and Selberg 2002].

