

Global Existence of Time-Dependent Yang-Mills-Higgs Monopoles

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Abstract. We study the Cauchy problem for non-abelian Yang-Mills-Higgs theory in $(3 + 1)$ -dimensional Minkowski spacetime. With suitable conditions on the background fields and a suitable choice of a Sobolev space for the subtracted gauge potentials and the Higgs field, we establish local existence. We then prove global existence by showing that an appropriate norm of the solutions cannot blow up in a finite time.

I. Introduction

Irving Segal [1], in 1963, introduced a general existence theory for semi-linear evolution equations. In 1979 Segal [2] himself showed that classical Yang-Mills theory could be cast into a suitable form to make use of this general theory and he showed that the local-in-time Cauchy problem could be solved. This means that if one is given regular initial data on \mathbb{R}^3 at some initial time (call it $t = 0$), there exists a unique smooth solution to the field equations compatible with the initial data over some finite time interval $(-t_0, t_0)$.

This result was improved on in 1982 by Ginibre and Velo [3] who added a Higgs field to the Yang-Mills potential and showed that local-in-time existence still held. The next major step forward was achieved, also in 1982, by Eardley and Moncrief [4] who independently derived the local existence result for Yang-Mills-Higgs theory and then extended this to obtain a global existence proof, that is to show that the time existence interval $(-t_0, t_0)$ can be made unboundedly large.

The major interest in classical solutions to the Yang-Mills-Higgs equations is due to the existence of magnetic monopole solutions (which demand that the Yang-Mills potential \mathbf{A} falls off like $1/r$) and non-trivial topologies in the Higgs field (which require that the Higgs field remain finite at infinity). However, all the existence results to date demand that both the Higgs field and the Yang-Mills potential be square-integrable on \mathbb{R}^3 , which naively demands that everything fall off faster than $r^{-3/2}$, and are incompatible with the magnetic monopoles and with non-trivial topologies.

In this paper we will give a global existence proof for a class of Yang-Mills-Higgs solutions which include solutions with magnetic monopoles and non-trivial topologies. The result we obtain here is not only interesting for its own sake; it is also immediately useful. For example, there has been some interest recently in approximation methods to study monopole-monopole scattering (e.g. Manton [5] and Atiyah and Hitchin [6]). In any approximation method one had better be sure that the approximate solution is an approximation to some exact solution, otherwise one may get nonsense (e.g. linearization instabilities in general relativity). The proof we give here should act as an underpinning to all these approximation techniques.

The method used in this paper to prove the desired existence theorem is quite straightforward. We specify, as part of the initial data a static background Yang-Mills potential \mathring{A} and Higgs field $\mathring{\Phi}$. These are chosen so that their asymptotic behaviour permits a finite magnetic charge and non-trivial topology. We then write the total potential A as $\mathring{A} + \mathbf{a}$ and Φ as $\mathring{\Phi} + \varphi$ and regard \mathbf{a} and φ as the dynamical fields which are square integrable, and so fall off rapidly at infinity. We will describe this in detail in Sect. II.

Of the two local-in-time existence techniques the Eardley and Moncrief method is superior to the Ginibre and Velo method. This is due to the fact that Eardley and Moncrief have found a way to build the $\text{div} E = 0$ constraint directly into the dynamical equations whereas Ginibre and Velo, following Segal, ignore it until the end, and then show that the dynamics are compatible with the constraint. This allows Eardley and Moncrief to prove existence with weaker conditions on the initial data than Ginibre and Velo require.

Unfortunately, the background subtraction method turns out to be incompatible with the Eardley-Moncrief local-in-time technique (which Eardley and Moncrief realised themselves.) Happily, the background subtraction method turns out to be entirely compatible with the Ginibre and Velo proof, and so our local existence proof is obtained by a straightforward extension of their technique. Of course, we have to include extra terms in our field equations which arise from the background static \mathring{A} , $\mathring{\Phi}$ fields, but these are easy to handle. We prove the local-in-time existence result in Sect. III.

To turn the local existence result into a global existence proof we switch back to Eardley and Moncrief and copy their technique. The difficulties in the local proof are not relevant to the global part and so we can follow them. Of course we still have to worry about the background fields. A further problem is that we have to make the global part of the proof agree with the local part. This means that we have to extend the global proof to one degree of differentiability higher than Eardley and Moncrief require. We manage to deal with both of these problems and get our global result in Sect. IV.

This should not come as any great surprise. For any theory, the only situations where one would have local but not global existence is where we would have either some local loss of differentiability through some form of singularity developing or where nothing bad would happen at any point, but some integral over the whole space would blow up. In view of the hyperbolic nature of the field equations, the singularity formation should be a strictly local phenomenon, and should not be influenced by asymptotic conditions. The problem of the norms blowing up can

only be avoided by sensible choice of norm. For us, the natural norms are the energy and the higher derivative pseudo-energies, and so we are not surprised when they remain finite.

II. Background and Dynamical Fields

We will study classical Yang-Mills-Higgs theory on Minkowski space with a Yang-Mills potential $A_\mu(t, \mathbf{x})$ and a Higgs scalar $\Phi(t, \mathbf{x})$. The field corresponding to A_μ is

$$F_{\mu\nu} = \partial_\mu A_\nu - \partial_\nu A_\mu + [A_\mu, A_\nu]. \quad (2.1)$$

The Lagrangian we will use is

$$\mathcal{L} = -\frac{1}{4}\langle F_{\mu\nu}, F^{\mu\nu} \rangle + \langle D_\mu \Phi, D^\mu \Phi \rangle - (|\Phi|^2 - 1)^2, \quad (2.2)$$

where D is the covariant derivative with A_μ as connection. This is the natural Lagrangian to use to get a non-trivial topology in the Higgs field at infinity and a magnetic monopole, together with finite energy. This is achieved by having Φ non-constant at infinity but in such a way that $|\Phi| \rightarrow 1$ at infinity. In addition we require $D_\mu \Phi = \partial_\mu \Phi + A_\mu \Phi$ to fall off rapidly at infinity. This can be achieved because while we naively expect $\partial_\mu \Phi$ to fall off like (at best) $1/r$ at infinity, we can assume $A_\mu \sim 1/r$ at infinity, and so $A_\mu \Phi$ also falls off like $1/r$. The two $1/r$ terms in $\partial_\mu \Phi$ and $A_\mu \Phi$ can cancel to give a faster fall-off to $D_\mu \Phi$. This is how the standard static monopole solutions behave.

We will assume the existence of two static background fields \hat{A}_μ and $\hat{\Phi}$. These, at least asymptotically, will behave like static monopole solutions. The actual potential and scalar fields we will write as

$$A_\mu(t, \mathbf{x}) = \hat{A}_\mu(\mathbf{x}) + a_\mu(t, \mathbf{x}), \quad \hat{A}_0 = 0, \quad (2.3a)$$

$$\Phi(t, \mathbf{x}) = \hat{\Phi}(\mathbf{x}) + \varphi(t, \mathbf{x}). \quad (2.3b)$$

We will treat the subtracted fields (a_μ, φ) as the dynamical fields and demand that they fall off rapidly enough to give us finite energy solutions.

We will describe everything in terms of the Sobolev spaces H_k with norm

$$\|u\|_{H_k}^2 = \sum_{|\alpha| \leq k} \|\partial^\alpha u\|_{L^2}^2 < \infty \quad (2.4)$$

for various values of k , except \hat{A} and $\hat{\Phi}$. We will assume that \hat{A} and $\hat{\Phi}$ satisfy the following conditions

- (i) $\hat{A}_i, \hat{\Phi} \in C^k$,
- (ii) $\hat{B}_i = \varepsilon_{ijk}(\partial_j \hat{A}_k + \hat{A}_j \hat{A}_k) \in H_{k+1}$,
 $\hat{\pi}_i = \partial_i \hat{\Phi} + \hat{A}_i \hat{\Phi} \in H_{k+1}$,
 $\partial_i \hat{A}_j \in H_{k-1}$, $(|\hat{\Phi}|^2 - 1) \in H_k$,

for some specified $k \geq 2$. These conditions are compatible with the magnetic charge associated with \hat{A}

$$\hat{g} = \oint_{\infty} \hat{\Phi} \cdot \hat{B}_i dx^i, \quad (2.5)$$

and the topological winding number associated with $\mathring{\Phi}$, i.e. for a triplet

$$\dot{n} = \frac{1}{8\pi} \oint_{\infty} \varepsilon_{ijk} \partial_j \mathring{\Phi}_{(a)} \partial_k \mathring{\Phi}_{(b)} \varepsilon_{abc} \mathring{\Phi}_{(c)} dx^i, \quad (2.6)$$

being non-zero. Although we are mainly interested in topologically nontrivial finite-energy configurations (magnetic monopoles), our proof covers topologically trivial configurations in topologically trivial and nontrivial models with spontaneous symmetry breaking as well.

III. Local Existence Proof

Following Segal we will work the local existence proof in the temporal gauge $A_0=0$. We will split $F_{\mu\nu}$ in the standard way,

$$E_i = F_{0i} = \partial_0 A_i, \quad (3.1)$$

$$B_i = \varepsilon_{ijk} (\partial_j A_k + A_j A_k). \quad (3.2)$$

We wish to make a further splitting, because we wish to write

$$A_i(t, \mathbf{x}) = \mathring{A}_i(\mathbf{x}) + a_i(t, \mathbf{x}), \quad (3.3)$$

where $\mathring{A}_i(\mathbf{x})$ is the static background with slow fall-off (with $\mathring{A}_0=0$), $a_i(t, \mathbf{x})$ is the time-varying field with fast fall-off. Having defined

$$\mathring{B}_i = \varepsilon_{ijk} (\partial_j \mathring{A}_k + \mathring{A}_j \mathring{A}_k), \quad (3.4)$$

we also define

$$b_i = B_i - \mathring{B}_i = \varepsilon_{ijk} (\partial_j a_k + a_j a_k + \mathring{A}_j a_k + a_j \mathring{A}_k), \quad (3.5)$$

and of course

$$e_i = E_i = \partial_0 a_i. \quad (3.6)$$

The natural splitting of the Higgs field follows from writing

$$\mathring{\Phi}(t, \mathbf{x}) = \mathring{\Phi}(\mathbf{x}) + \varphi(t, \mathbf{x}). \quad (3.7)$$

In general we have for the first derivative of $\mathring{\Phi}$,

$$\pi = D_0 \mathring{\Phi} = \partial_0 \mathring{\Phi}, \quad (3.8)$$

$$\pi_i = D_i \mathring{\Phi} = \partial_i \mathring{\Phi} + A_i \mathring{\Phi}. \quad (3.9)$$

We now define

$$\psi = \partial_0 \varphi = \pi, \quad (3.10)$$

$$\psi_i = \pi_i - \mathring{\pi}_i = \partial_i \varphi + a_i \varphi + a_i \mathring{\Phi} + \mathring{A}_i \varphi. \quad (3.11)$$

The initial data (at a fixed time $t=0$) we will specify for the Yang-Mills-Higgs field will consist of the static background fields $(\mathring{A}_i, \mathring{\Phi})$ and a sextuplet of fields

$$(a_i, e_i, b_i, \varphi, \psi, \psi_i) = u(t, \mathbf{x}). \quad (3.12)$$

Following Ginibre and Velo [3] [their Eqs. (2.4) and (2.5)], the dynamical equations now read

$$\partial_0 u(t) = Tu(t) + K(u(t)) + \mathring{K}(u(t)), \quad (3.13)$$

with

$$T = \begin{pmatrix} T_a & 0 \\ 0 & T_\varphi \end{pmatrix}, \quad (3.14)$$

$$T_a = \begin{pmatrix} 0 & 1 & 0 \\ 0 & 0 & R_a \\ 0 & -R_a^* & 0 \end{pmatrix}, \quad T_\varphi = \begin{pmatrix} 0 & 1 & 0 \\ 0 & 0 & R \\ 0 & -R^* & 0 \end{pmatrix},$$

$$(R_a)_{ij} = \varepsilon_{ijk} \partial_k, \quad R = -V, \quad (3.15)$$

$$K = (0, -\varepsilon_{ijk}[a_j, b_k] - 2 \operatorname{Re} \langle \psi_i, \theta_a \varphi \rangle \theta_a, \varepsilon_{ijk}(a_j e_k + a_j a_k), 0, a_i \psi_i - |\varphi|^2 \varphi, e_i \varphi + a_i \psi), \quad (3.16)$$

$$\begin{aligned} \mathring{K} = & (0, -\varepsilon_{ijk}(\partial_j \mathring{B}_k + [\mathring{A}_j, b_k] + [a_j, \mathring{B}_k] + [\mathring{A}_j, \mathring{B}_k]) \\ & - 2 \operatorname{Re} \langle \psi_i, \theta_a \mathring{\Phi} \rangle + \langle \mathring{\pi}_i, \theta_a \varphi \rangle \\ & + \langle \mathring{\pi}_i, \theta_a \mathring{\Phi} \rangle) \theta_a, \varepsilon_{ijk}(\mathring{A}_j e_k + e_j \mathring{A}_k), \\ & 0, \partial_i \mathring{\pi}_i + \mathring{A}_i \psi_i + a_i \mathring{\pi}_i - 2 \mathring{\Phi} |\varphi|^2 \\ & - 2(\varphi + \mathring{\Phi})(|\mathring{\Phi}|^2 + 2 \langle \mathring{\Phi}, \varphi \rangle + |\varphi|^2 - 1), e_i \mathring{\Phi} + \mathring{A}_i \psi). \end{aligned} \quad (3.17)$$

These, by themselves, do not constitute the Yang-Mills-Higgs field equations. We must add three constraints to them

$$b_i - \varepsilon_{ijk}(\partial_j a_k + a_j a_k + \mathring{A}_j a_k + a_j \mathring{A}_k) = 0, \quad (3.18)$$

$$\psi_i - \partial_i \varphi - a_i \varphi - a_i \mathring{\Phi} - \mathring{A}_i \varphi = 0, \quad (3.19)$$

$$\begin{aligned} \partial_i e_i + [a_i, e_i] + 2 \operatorname{Re} \langle \psi, \theta_a \varphi \rangle \theta_a \\ + [\mathring{A}_i, e_i] + 2 \operatorname{Re} \langle \psi, \theta_a \mathring{\Phi} \rangle \theta_a = 0. \end{aligned} \quad (3.20)$$

These equations are essentially equivalent to those of Ginibre and Velo. The only difference is that their variables are

$$(A_i, E_i, B_i, \Phi, \pi, \pi_i) \quad \text{not} \quad (a_i, e_i, b_i, \varphi, \psi, \psi_i).$$

With the change due to subtracting the background, we still have that the linear operator T is identical to theirs, and our K operator is identical to their non-linear term. The only change is the extra \mathring{K} term, where we accumulate all terms which depend on the background.

We wish to apply the Segal existence theorem to (3.13). This theorem will not deal directly with (3.13), but with the associated integral equation

$$u(t) = U(t)u_0 + \int_0^t d\tau U(t-\tau)K_{\text{tot}}(u(\tau)), \quad (3.21)$$

where

$$K_{\text{tot}} = K + \mathring{K}, \quad (3.22)$$

and where U is the transformation generated by the linear operator T [(3.14) and (3.15)].

The first thing we have to do is to choose a Banach space B in which we want $u(t)$ to lie. In this case the obvious choice is

$$B_k = H_k(\mathbb{R}^3) \times H_k \times \dots \times H_k, \quad (3.23)$$

where k is some positive integer. We now show that U is very well behaved on B .

Lemma 2.1 of Ginibre and Velo give us the desired result

Lemma 3.1. *For any k , $U(t)$ is a (bounded) strongly continuous one parameter group in B_k and for any $t \in \mathbb{R}$, $u_0 \in B_k$, $U(t)$ satisfies the estimate*

$$\begin{aligned} \|U(t)u_0\| &\leq \mu(t) \|u_0\|, \\ \mu(t) &= \{1 + \frac{1}{2}|t|(|t| + (t^2 + 4)^{1/2})\}^{1/2}. \end{aligned}$$

We next have to look at K_{tot} and show that K_{tot} maps B into itself. To do this we choose $k \geq 2$. This allows us to use the fact that $H_k(\mathbb{R}^3)$ is a Banach algebra for $k \geq 2$. There are no derivatives in K , only terms like $a_i e_j$. Since $a_j \in H_k$ and $e_j \in H_k \Rightarrow a_i e_j \in H_k$, obviously K maps $B \rightarrow B$. In \hat{K} there are two derivative terms $\partial_j \hat{B}$ and $\partial_i \hat{\pi}$. We put as part of the initial conditions that $\hat{B}, \hat{\pi} \in H_{k+1}$. The other terms are of the kind $\hat{A}_i e_j$. If we demand $\hat{A}_i \in C^k$, then $e_j \in H_k \Rightarrow \hat{A}_i e_j \in H_k$. The only other term that needs special treatment is $\hat{\Phi}^2 - 1 \in H_k$. It then follows that the operator K_{tot} maps B into itself.

The second property that we need for K_{tot} is that it satisfy a Lipschitz condition

$$\|K_{\text{tot}}(u_1) - K_{\text{tot}}(u_2)\| \leq C(\|u_1\|, \|u_2\|) \|u_1 - u_2\|, \quad (3.24)$$

where C is an increasing finite real function, and $\|\cdot\|$ is the $(H_k)^6$ norm. This again follows from the Banach Algebra property of H_k for $k \geq 2$. All the terms in K are of the form ab except for the φ^3 term. The ab term leads to

$$a_1 b_1 - a_2 b_2 = a_1(b_1 - b_2) - b_2(a_1 - a_2)$$

and

$$\|a_1(b_1 - b_2)\|_{H_k} \leq C_0 \|a_1\|_{H_k} \|b_1 - b_2\|_{H_k}.$$

The φ^3 term gives

$$\|\varphi_1^3 - \varphi_2^3\| = (\varphi_1 - \varphi_2)(\varphi_1^2 + \varphi_1\varphi_2 + \varphi_2^2),$$

and

$$\begin{aligned} \|(\varphi_1^2 + \varphi_1\varphi_2 + \varphi_2^2)(\varphi_1 - \varphi_2)\|_{H_k} &\leq C_0 \|\varphi_1^2 + \varphi_1\varphi_2 + \varphi_2^2\|_{H_k} \|\varphi_1 - \varphi_2\|_{H_k} \\ &\leq C_1 (\|\varphi_1\| + \|\varphi_2\|)^2 \|\varphi_1 - \varphi_2\|_{H_k}. \end{aligned}$$

The terms in \hat{K} are of the form $\hat{A}b$ which gives $\hat{A}(b_1 - b_2)$ and

$$\|\hat{A}(b_1 - b_2)\|_{H_k} \leq \|\hat{A}\|_{C^k} \|b_1 - b_2\|_{H_k},$$

together with $\hat{\Phi}^2$ which gives $\hat{\Phi}(\varphi_1^2 - \varphi_2^2)$ and satisfies

$$\|\hat{\Phi}(\varphi_1^2 - \varphi_2^2)\|_{H_k} \leq C_2 \|\hat{\Phi}\|_{C^k} (\|\varphi_1\| + \|\varphi_2\|) \|\varphi_1 - \varphi_2\|_{H_k}.$$

All these together combine to prove (3.24).

The Lipschitz condition on K_{tot} and the fact that $U(t)$ is a (semi)group is all that we need to apply the Segal existence theory. This immediately gives us that there exists a unique solution to (3.21) for some finite time interval $[-T, T]$ and that $u(t)$ for any $t \in [-T, T]$ belongs to H_k .

To go from the integral Eq. (3.21) to the differential equation (3.13), we have to show that K_{tot} is a C^1 map. This is quite straightforward. Since K_{tot} is a third order polynomial in u (the φ^3 term), the Frechet derivative of K_{tot} is quadratic in u . We now just repeat our proof of the Lipschitz condition for this simpler case. Analogously, we can prove that K_{tot} is a C^∞ map since all derivatives higher than third order vanish.

As we remarked earlier, the differential Eq. (3.13) is not equivalent to the Yang-Mills-Higgs equations. We need to also satisfy the constraints (3.18), (3.19), (3.20). This is quite straightforward. We choose the initial data at $t=0$ so as to satisfy the constraints. The dynamical Eq. (3.13) is consistent with the constraints and automatically propagates them (see Proposition 2.3 of Ginibre and Velo). Thus we have our desired local existence proof.

IV. A Priori Bounds

To expand the local proof to a global existence proof we will switch over to imitate the global existence proof of Eardley and Moncrief [4]. The key idea in their analysis is to obtain an a priori bound on $F_{\mu\nu}$ and $D_\mu\Phi$.

We start off with regular initial data and use the local existence proof to demonstrate the existence of a classical solution on a patch of spacetime. (To guarantee the smoothness Eardley and Moncrief want for working in Cronström gauge we impose the necessary smoothness conditions on the initial value data.) To extend this to a global solution we have to show that nothing goes bad at the boundary. This requires showing that the H_2 norms of $(a_i, e_i, b_i, \varphi, \psi, \psi_i)$ do not blow up. This means that we extend beyond the boundary and so continue.

Just as with the local Ginibre and Velo proof [3], we need only make minor modifications to the global Eardley and Moncrief proof to accommodate the extra \dot{A} and $\dot{\Phi}$ terms. The first point of importance is that we still have a finite conserved

energy $E_0 = \int_{\mathbb{R}^3} \varepsilon d^3x$, where

$$\varepsilon = \frac{1}{2}(E^2 + B^2) + \pi \cdot \pi + D_i\Phi \cdot D^i\Phi + (|\Phi|^2 - 1)^2 > 0. \quad (4.1)$$

This conserved energy can be derived from an energy-momentum tensor $T^{\mu\nu}$ given by

$$T^{\mu\nu} = \text{Tr} \{ F^{\mu\alpha} F_\alpha^\nu - \frac{1}{4} \eta^{\mu\nu} F_{\alpha\beta} F^{\alpha\beta} \} + 2D^\mu\Phi \cdot D^\nu\Phi - \eta^{\mu\nu} D_\alpha\Phi \cdot D^\alpha\Phi - \eta^{\mu\nu} (|\Phi|^2 - 1)^2, \quad (4.2)$$

which satisfies $\partial_\nu T^{\mu\nu} = 0$.

Let us choose a point p in the domain of local existence and shift coordinates so that p becomes the origin of coordinates. The initial slice now is labeled $t = -t_0$. Consider the back light cone K_p from p to the original surface. Call the solid sphere where the light-cone intersects the original surface B_p . Now consider the

conservation of energy on the set which is the interior of K_p (call it \tilde{K}_p)

$$0 = \int_{\tilde{K}_p} \partial_\nu T_0^\nu d^4x = \int_{K_p} T_0^\nu \cdot n + \int_{B_p} T_0^\nu \cdot n \quad (4.3)$$

(on using the Gauss theorem). This, when written out, is of the form

$$\begin{aligned} & \int_{K_p} r^2 dr d\Omega \left\{ \frac{1}{2} \text{Tr} \left[\frac{3}{4} (\hat{\ell} \cdot F \cdot \hat{m})^2 + (\hat{\ell} \cdot F \cdot \hat{e}_A)^2 \right. \right. \\ & \quad \left. \left. + \frac{1}{2} (\hat{e}_A \cdot F \cdot \hat{e}_B)^2 \right] + (|\Phi|^2 - 1)^2 \right. \\ & \quad \left. + [D_{\hat{e}} \Phi \cdot D_{\hat{e}} \Phi + D_{\hat{e}_A} \Phi \cdot D_{\hat{e}_A} \Phi] \right\} \\ & = \int_{B_p} \varepsilon d^3x < E_0, \end{aligned} \quad (4.4)$$

where

$$\begin{aligned} \hat{\ell} &= -\partial/\partial t + \partial/\partial r, & \hat{m} &= \partial/\partial t + \partial/\partial r, \\ \hat{e} &= \left(\frac{1}{r} \partial/\partial \vartheta, \frac{1}{r \sin \vartheta} \partial/\partial \varphi \right). \end{aligned}$$

The integral along the light-cone is a set of quadratic (positive) terms. Therefore we have that each one individually is bounded by the finite total energy E_0 , i.e.

$$\begin{aligned} & \int_{K_p} r^2 dr d\Omega \text{Tr}(\hat{\ell} \cdot F \cdot \hat{m})^2 < E_0, \\ & \int_{K_p} r^2 dr d\Omega (|\Phi|^2 - 1)^2 < E_0, \end{aligned} \quad (4.5)$$

and so on.

The next stage is to realise that the Yang-Mills equations can be manipulated to give us wave equations for F and $D\Phi$:

$$\begin{aligned} \eta^{\mu\nu} \partial_\mu \partial_\nu F_{\alpha\beta} &= -2\partial_\gamma [A^\gamma, F_{\alpha\beta}] + [\partial_\gamma A^\gamma, F_{\alpha\beta}] \\ & \quad - [A^\gamma, [A_\gamma, F_{\alpha\beta}]] + 2[F_{\alpha\gamma}^\gamma, F_{\gamma\beta}] + 2((F_{\alpha\beta} \Phi) \cdot \theta_\alpha \Phi) \theta_\alpha \\ & \quad + 2((D_\beta \Phi) \cdot \theta_\alpha (D_\alpha \Phi) - 2(D_\alpha \Phi) \cdot \theta_\alpha (D_\beta \Phi)) \theta_\alpha = : \varrho_{\alpha\beta}, \end{aligned} \quad (4.6)$$

$$\begin{aligned} \eta^{\mu\nu} \partial_\mu \partial_\nu (D_\alpha \Phi) &= -2\partial_\mu (A^\mu D_\alpha \Phi) + (\partial_\mu A^\mu) D_\alpha \Phi - A_\mu A^\mu D_\alpha \Phi \\ & \quad + 2((D_\alpha \Phi) \cdot \theta_\alpha \Phi) \theta_\alpha \Phi - 2F_\alpha^\mu D_\mu \Phi + 4D_\alpha \{ \Phi (|\Phi|^2 - 1) \} = : \varrho_\alpha. \end{aligned} \quad (4.7)$$

The right hand side of each of these equations can be regarded as a source term for a wave equation. This allows us to write $F_{\alpha\beta}$ at p as the average value of $F_{\alpha\beta}$ on the boundary of B_p plus a term arising from the integral of the source along K_p . In other words

$$F_{\alpha\beta}(p) = \frac{1}{4\pi} \int_{S^2} d\Omega (r_0 m^\mu \partial_\mu F_{\alpha\beta} + F_{\alpha\beta}) - \frac{1}{4\pi} \int_{K_p} r dr d\Omega \varrho_{\alpha\beta}, \quad (4.8)$$

where S_2 is the sphere of radius $r = r_0 = -t_0$ on the original slice. Equivalently we have

$$D_\mu \Phi(p) = \frac{1}{4\pi} \int_{S^2} d\Omega (r_0 m^2 \partial_\alpha D_\mu \Phi + D_\mu \Phi) - \frac{1}{4\pi} \int_{K_p} r dr d\Omega \varrho_\mu. \quad (4.9)$$

There is no difficulty in handling the two-sphere integral and showing that it is finite. The real problem arises from the light-cone integral, and there we need to use the bounds obtained from the energy conservation equation. One useful bound that can be obtained is

$$\begin{aligned} \int_{K_p} (|\Phi|^2 - 1)^2 < E_0 &\Rightarrow \int_{K_p} |(|\Phi|^2 - 1)| < \sqrt{CE_0} t_0^{3/2} \\ &\Rightarrow \int_{K_p} |\Phi|^2 < \sqrt{CE_0} t_0^{3/2} + Ct_0^3 \end{aligned} \quad (4.10)$$

(where Ct_0^3 is the volume of the cone) and also

$$\int_{K_p} |\Phi|^4 < E_0 + 2\sqrt{CE_0} t_0^{3/2} + 3Ct_0^3. \quad (4.11)$$

The only real problem we have in going from Eardley-Moncrief to our spaces is the fact that for us such functions as

$$\int_{\mathbb{R}^3} |\Phi|^2 d^3x$$

are not finite.

However, in deriving the a priori bounds on F and $D\Phi$ we need never integrate over \mathbb{R}^3 only over K_p , where the relevant integrals are finite and so the proof goes through unchanged. Thus we also have that $\|F\|_{L^\infty}$ and $\|D_\mu\Phi\|_{L^\infty}$ remain finite.

V. Bounds on Everything Else

Given that $\|F\|_{L^\infty}$ and $\|D_\mu\Phi\|_{L^\infty}$ remain finite in the local range of existence of a solution, we now need to show that the H_2 norms of $(a_i, e_i, b_i, \varphi, \psi, \psi_i)$ also remain finite. We begin by showing that the L_2 norms remain finite. Let us begin with

$$\frac{d}{dt} (\|\varphi\|_{L_2})^2 = 2 \int_{\mathbb{R}^3} \psi \cdot \varphi \leq 2 \left| \int_{\mathbb{R}^3} \psi^2 \right|^{1/2} \left| \int_{\mathbb{R}^3} \varphi^2 \right|^{1/2} \leq \sqrt{2E_0} \|\varphi\|_{L_2}, \quad (5.1)$$

which implies

$$\|\varphi(t)\|_{L_2} \leq \|\varphi(0)\|_{L_2} + \sqrt{2E_0} t, \quad (5.2)$$

and so $\|\varphi\|_{L_2}$ remains finite. Next

$$\frac{d}{dt} (\|a_i\|_{L_2})^2 = 2 \int_{\mathbb{R}^3} a_i \cdot e_i \leq \sqrt{2E_0} \|a_i\|_{L_2}. \quad (5.3)$$

The L_2 norms of $e_i, \psi, B_i = b_i + \mathring{B}_i$ and $D_i\Phi = \psi_i + \mathring{\pi}_i$ do not blow up from energy conservation. Because \mathring{B}_i and $\mathring{\pi}_i$ are constant and belong to L_2 , the L_2 norms of b_i and ψ_i cannot blow up.

The next stage is to integrate the equations

$$\frac{\partial}{\partial t} a_i = e_i, \quad \frac{\partial}{\partial t} \varphi = \psi \quad (5.4)$$

to give

$$\|a_i\|_{L^\infty} \leq \|a_i(0)\|_{L^\infty} + \int_0^t ds \|e_i(s)\|_{L^\infty} < \infty, \quad (5.5a)$$

and

$$\|\varphi\|_{L^\infty} \leq \|\varphi(0)\|_{L^\infty} + \int_0^t ds \|\psi(s)\|_{L^\infty}. \quad (5.5b)$$

The final inequality follows from the fact that $e_i(s)$ is part of $F_{\alpha\beta}$ and $\psi(s)$ is part of $D\Phi$, both of which were shown to remain finite in Sect. IV.

We have that

$$b_i = B_i - \dot{B}_i = \varepsilon_{ijk}(\partial_j a_k + a_j a_k + \dot{A}_j a_k + a_j \dot{A}_k), \quad (5.6a)$$

and

$$\psi_i = \pi_i - \dot{\pi}_i = \partial_i \varphi + a_i \varphi + a_i \dot{\Phi} + \dot{A}_i \varphi. \quad (5.6b)$$

We know $a_i, b_i, \psi_i, \varphi \in L_2$, and $\|a_i\|_{L^\infty}, \|\dot{A}_i\|_{L^\infty}$ and $\|\dot{\Phi}\|_{L^\infty}$ are finite. Therefore, $\partial_i \varphi \in L_2$ and $\varepsilon_{ijk} \partial_j a_k \in L_2$ hold. Further $\|\psi_i\|_{L^\infty}, \|\partial_i \varphi\|_{L^\infty}, \|b_i\|_{L^\infty}$, and $\|\varepsilon_{ijk} \partial_j a_k\|_{L^\infty}$ are all finite. Finally the third constraint

$$\partial_i e_i = -[a_i, e_i] - 2 \operatorname{Re} \langle \psi, \theta_a \varphi \rangle \theta_a - [\dot{A}_i, e_i] - 2 \operatorname{Re} \langle \psi, \theta_a \dot{\Phi} \rangle \theta_a \quad (5.7)$$

also gives us $\partial_i e_i \in L_2$, and $\|\partial_i e_i\|_{L^\infty}$ finite.

The next stage is to show that $(a_i, e_i, b_i, \varphi, \psi, \psi_i) \in H_1$. This means

$$\partial_i a_j, \partial_i e_j, \partial_i b_j, \partial_i \varphi, \partial_i \psi, \partial_i \psi_j \in L_2.$$

Let us start with $\partial_i a_j$ and consider

$$\mathcal{E}_0 = -\frac{1}{2} \operatorname{Tr} \int_{\mathbb{R}^3} [e_i \cdot e_i + (\partial_i a_j)(\partial_i a_j)] d^3 x. \quad (5.8)$$

Now

$$\begin{aligned} \frac{d}{dt} \mathcal{E}_0 &= -\operatorname{Tr} \int [e_i \{ \varepsilon_{ijk} \partial_k b_i - \varepsilon_{ijk} [a_j, b_k] \\ &\quad - 2 \operatorname{Re} \langle \psi_i, \theta_a \varphi \rangle \theta_a + \varepsilon_{ijk} (\partial_k \dot{B}_j - [\dot{A}_j, b_k] \\ &\quad - [a_j, \dot{B}_k] - [\dot{A}_j, \dot{B}_k] - 2 \operatorname{Re} \langle \psi_i, \theta_a \varphi \rangle \\ &\quad + \langle \dot{\pi}_i, \theta_a \varphi \rangle + \langle \dot{\pi}_i, \theta_a \dot{\Phi} \rangle \theta_a \} + \partial_i a_j \partial_i e_j] d^3 x. \end{aligned} \quad (5.9)$$

All the terms in the expression above involving background fields are clearly finite and the two terms $e_i \varepsilon_{ijk} [a_j, b_k]$ and $e_i \operatorname{Re} \langle \psi_i, \theta_a \varphi \rangle \theta_a$ are finite because $\|e_i\|_{L^\infty}, \|a_j\|_{L_2}, \|b_k\|_{L_2}, \|\psi_i\|_{L_2}, \|\varphi\|_{L_2} < \infty$.

This leaves only

$$\begin{aligned} \int \operatorname{Tr} \{ e_i \varepsilon_{ijk} \partial_k b_i + \partial_i a_j \partial_i e_j \} &= \int \operatorname{Tr} \{ e_i \partial_k (\partial_k a_i - \partial_i a_k + [a_k, a_i] \\ &\quad + [\dot{A}_k, a_i] + [a_k, \dot{A}_i]) + \partial_i a_j \partial_i e_j \}. \end{aligned} \quad (5.10)$$

The first and last terms combine to form $\partial_k (e_i \partial_k a_i)$ and so integrate to zero. The second term can be written as

$$\partial_i e_i \partial_k a_k - \partial_i (e_i \partial_i a_k).$$

Now

$$\int \partial_i e_i \partial_k a_k \leq \| \partial_i e_i \|_{L_2} \| \partial_k a_k \|_{L_2} \leq C \sqrt{\mathcal{E}_0}, \quad (5.11)$$

since $\| \partial_i e_i \|_{L_2}$ is finite. The total divergence can again be ignored. The three commutator terms fall in line, the worst being

$$\int e_i a_k \partial_k a_i \leq \| e_i \|_{L^\infty} \| a_k \|_{L_2} \| \partial_k a_i \|_{L_2} \leq C_1 \sqrt{\mathcal{E}_0}. \quad (5.12)$$

Hence we have

$$\frac{d}{dt} \mathcal{E}_0 = C_2 + C_3 \sqrt{\mathcal{E}_0}, \quad (5.13)$$

which implies that \mathcal{E}_0 remains finite and hence $\partial_i a_k \in L_2$.

To show that e_i , b_i , ψ , and ψ_i belong to H_1 we have to show that

$$\mathcal{E}_1 = \frac{1}{2} \int \{ (\partial_j e_i) (\partial_j e_i) + (\partial_j b_i) (\partial_j b_i) + (\partial_i \psi) (\partial_i \psi) + (\partial_i \psi_j) (\partial_i \psi_j) \} d^3 x \quad (5.14)$$

remains finite. To show this we calculate

$$\begin{aligned} \frac{d}{dt} \mathcal{E}_1 = & \int \{ \partial_j e_i \partial_j [\varepsilon_{imn} \partial_n b_m - \varepsilon_{imn} [a_m, b_n] \\ & - 2 \operatorname{Re} \langle \psi_i, \theta_a \varphi \rangle \theta_a - \varepsilon_{imn} (\partial_m \dot{B}_n + [a_m, \dot{B}_n] \\ & + [\dot{A}_m, b_n] + [\dot{A}_m, \dot{B}_n]) - 2 \operatorname{Re} \langle \psi_i, \theta_a \dot{\Phi} \rangle \\ & + \langle \dot{\pi}_i, \theta_a \varphi \rangle + \langle \dot{\pi}_i, \theta_a \dot{\Phi} \rangle \theta_a \\ & + \partial_j b_i \partial_j [\varepsilon_{imn} \partial_m e_n + \varepsilon_{imn} [a_m, e_n] \\ & + \varepsilon_{imn} [\dot{A}_m, e_n]] + (\psi / \psi_i - \text{terms}) \}. \end{aligned} \quad (5.15)$$

The two leading terms form a total divergence $\partial_n \{ \partial_j e_i \varepsilon_{imn} \partial_j b_m \}$ and so can be ignored. The next terms to be considered are

$$\int \partial_j e_i \partial_j a_m b_n \leq \| b_n \|_{L^\infty} \| \partial_j e_i \|_{L_2} \| \partial_j a_m \|_{L_2} \leq C \sqrt{\mathcal{E}_1}, \quad (5.16)$$

and

$$\int a_m \partial_j e_i \partial_j b_n \leq \| a_m \|_{L^\infty} \| \partial_j e_i \|_{L_2} \| \partial_j b_n \|_{L_2} \leq C \mathcal{E}_1. \quad (5.17)$$

The terms $\partial_j e_i \partial_j \psi_i \varphi$ and $\psi_i \partial_j e_i \partial_j \varphi$ can be handled in the same way.

All the terms involving the background fields are quite straightforward. We only need to use

$$\partial_i \partial_j \dot{B}_m, \partial_j \dot{A}_m \in L_2; \quad \dot{A}_m, \dot{B}_m \in L^\infty$$

to get terms proportional to $\sqrt{\mathcal{E}_1}$ and to \mathcal{E}_1 . We finally get

$$\frac{d}{dt} \mathcal{E}_1 \leq C_1 \sqrt{\mathcal{E}_1} + C_2 \mathcal{E}_1, \quad (5.18)$$

which shows that \mathcal{E}_1 remains finite, and hence

$$\| \partial_j e_i \|_{L_2}, \| \partial_j b_i \|_{L_2}, \| \partial_i \psi \|_{L_2}, \| \partial_i \psi_j \|_{L_2}$$

remain finite. These can now be inserted into the constraints to force $\varepsilon_{ijk}\partial_j a_k$, $\partial_i \varphi$, and $\partial_i e_i$ to belong to H_1 .

The next step is to show that $\partial_j a_k$ itself belongs to H_1 , or equivalently that $\partial_i \partial_j a_k$ belongs to L_2 . This is the same as showing that

$$\mathcal{E}'_1 = \frac{1}{2} \int \{(\partial_j e_i)(\partial_j e_i) + (\partial_i \partial_j a_k)(\partial_i \partial_j a_k)\} \quad (5.19)$$

remains finite. To do this we calculate

$$\frac{d}{dt} \mathcal{E}'_1 = \int \{(\partial_j \dot{e}_i)(\partial_j \partial_0 e_i) + (\partial_i \partial_j \dot{a}_k)(\partial_i \partial_j e_k)\}. \quad (5.20)$$

Substituting for $\partial_0 e_i$ we get a total derivative plus

$$\partial_j(\partial_i e_i)\partial_j \partial_n a_n - (\partial_j e_i)\partial_j \partial_n [a_i, a_n] + \dots$$

On using $\partial_i e_i \in H_1$ we can bound the first of these terms by $C\sqrt{\mathcal{E}'_1}$.

The second term requires some care. One part will be of the form

$$\int a_i \partial_j e_i \partial_j \partial_n a_n \leq \|a_i\|_{L^\infty} \|\partial_j e_i\|_{L_2} \|\partial_j \partial_n a_n\|_{L_2} \leq C\mathcal{E}'_1. \quad (5.21)$$

The other parts, however, are of the form

$$\partial_j e_i \partial_j a_i \partial_n a_n \quad \text{or} \quad \partial_j a_i \partial_n a_i \partial_j a_n.$$

The trick is to remember that $\|F_{\mu\nu}\|_{L^\infty}$ is finite which implies $\|b_i\|_{L^\infty}$ finite which implies $\|\partial_j a_i - \partial_i a_j\|_{L^\infty}$ finite, and so we can replace $\partial_j a_i$ with $\partial_i a_j$ to give

$$\partial_j e_i \partial_i a_j \partial_n a_n = \partial_i [\partial_j e_i a_j \partial_n a_n] - a_j \partial_n a_n \partial_j \partial_i e_i - a_j \partial_j e_i \partial_i \partial_n a_n. \quad (5.22)$$

The divergence is fine, the first term is finite because $\partial_i e_i \in H_1$ and the second term is proportional to $\sqrt{\mathcal{E}'_1}$. Everything else is well behaved, so we finally get

$$(a_i, \varphi) \in H_2, \quad (e_i, b_i, \psi, \psi_i) \in H_1.$$

The last stage of the proof is to show that the $H_2 \times H_2$ norm remains finite. The obvious thing to do is to write down the equivalent of \mathcal{E}_1 , i.e.

$$\begin{aligned} \mathcal{E}_2 = \frac{1}{2} \int \{ & (\partial_k \partial_j e_i)(\partial_k \partial_j e_i) + (\partial_k \partial_j b_i)(\partial_k \partial_j b_i) + (\partial_k \partial_j \psi)(\partial_k \partial_j \psi) \\ & + (\partial_k \partial_j \psi_i)(\partial_k \partial_j \psi_i) \} d^3x \end{aligned} \quad (5.23)$$

and show $d\mathcal{E}_2/dt \leq$ (polynomial with constant coefficients in \mathcal{E}_2). Now

$$\begin{aligned} \frac{d}{dt} \mathcal{E}_2 = \int & \partial_k \partial_j \dot{e}_i \partial_k \partial_j [\varepsilon_{imn} \partial_n b_m - \varepsilon_{imn} [a_m, b_n]] \\ & - 2 \operatorname{Re} \langle \psi_i, \theta_a \varphi \rangle \theta_a - \varepsilon_{imn} (\partial_m \dot{B}_n + [a_m, \dot{B}_n] + [\dot{A}_m, b_n] + [\dot{A}_m, \dot{B}_n]) \\ & - 2 \operatorname{Re} \langle \psi_i, \theta_a \dot{\Phi} \rangle + \langle \dot{\pi}_i, \theta_a \varphi \rangle + \langle \dot{\pi}_i, \theta_a \dot{\Phi} \rangle \theta_a \\ & + \partial_k \partial_j \dot{b}_i \partial_j \partial_k [\varepsilon_{imn} \partial_m e_n + \varepsilon_{imn} [a_m, e_n] + \varepsilon_{imn} [\dot{A}_m, e_n]] \\ & + (\psi/\psi_i - \text{terms}). \end{aligned} \quad (5.24)$$

Again, the leading terms form total divergences

$$\partial_n[\varepsilon_{imn}\partial_k\partial_j e_i\partial_k\partial_j b_m] \quad \text{and} \quad \partial_i[\partial_k\partial_j\psi\partial_k\partial_j\psi_i],$$

and so can be ignored. However when we look at the next term we run into difficulties. This is of the form

$$\partial_k\partial_j e_i \varepsilon_{imn}\partial_k\partial_j [a_m, b_n].$$

The two terms

$$\int \varepsilon_{imn} b_n \partial_k \partial_j e_i \partial_k \partial_j a_m \leq \|b_n\|_{L^\infty} \|\partial_k \partial_j e_i\|_{L_2} \|\partial_k \partial_j a_m\|_{L_2} \leq C\sqrt{\mathcal{E}_2}, \quad (5.25)$$

and

$$\int \varepsilon_{imn} a_m \partial_k \partial_j e_i \partial_k \partial_j b_n \leq \|a_m\|_{L^\infty} \|\partial_k \partial_j e_i\|_{L_2} \|\partial_k \partial_j b_n\|_{L_2} \leq C\mathcal{E}_2 \quad (5.26)$$

offer no difficulty, but the third term $\int \varepsilon_{imn} \partial_k \partial_j e_i \partial_j a_m \partial_k b_n$ has to be dealt with carefully. We know that it is bounded by

$$\|\partial_k \partial_j e_i\|_{L_2} \|\partial_j a_m \partial_k b_n\|_{L_2}.$$

We also know from the Sobolev multiplication theorem that $H_1 \times H_1$ is uniformly embedded in L_2 . Hence

$$\|\partial_j a_m \partial_k b_n\|_{L_2} \leq C_0 \|\partial_j a_m\|_{H_1} \|\partial_k b_n\|_{H_1}. \quad (5.27)$$

But of course $\|\partial_j a_m\|_{H_1}$ is finite and $\|\partial_k b_n\|_{H_1}$ is bounded by a finite number times $\sqrt{\mathcal{E}_2}$. Hence

$$\int \varepsilon_{imn} \partial_k \partial_j e_i \partial_j a_m \partial_k b_n \leq C_1 \sqrt{\mathcal{E}_2} + C_2 \mathcal{E}_2. \quad (5.28)$$

All the other terms are well behaved. To deal with the terms which include the background terms we need slightly sharper conditions, on these, i.e.

$$\mathring{B} \in H_3, \quad \mathring{\pi}_i \in H_3, \quad \partial_j \mathring{A} \in L_2, \quad (5.29)$$

and we finally get

$$\frac{d}{dt} \mathcal{E}_2 \leq C_1 + C_2 \sqrt{\mathcal{E}_2} + C_3 \mathcal{E}_2. \quad (5.30)$$

This guarantees that \mathcal{E}_2 remains finite and hence that each of (e_i, b_i, ψ, ψ_i) remain in H_2 if originally there. We already know a_i remains in H_2 , and of course since $\psi_i = \partial_i \varphi + a_i \varphi + a_i \mathring{\Phi} + \mathring{A}_i \varphi$, we know that φ belongs to H_3 . Hence we have that the norms we needed to prove the local existence theorem do not blow up, and this is sufficient to prove global existence.

Acknowledgement. JB would like to thank the Alexander-von-Humboldt Foundation for the partial financial support he receives as a Feodor Lynen Fellow.

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Communicated by L. Nirenberg

Received October 29, 1985