NON-CMC SOLUTIONS TO THE EINSTEIN CONSTRAINT EQUATIONS ON ASYMPTOTICALLY EUCLIDEAN MANIFOLDS WITH APPARENT HORIZON BOUNDARIES

MICHAEL HOLST AND CALEB MEIER

ABSTRACT. In this article we further develop the solution theory for the Einstein constraint equations on an n-dimensional, asymptotically Euclidean manifold \mathcal{M} with interior boundary Σ . Building on recent results for both the asymptotically Euclidean and compact with boundary settings, we show existence of far-from-CMC and near-CMC solutions to the conformal formulation of the Einstein constraints when nonlinear Robin boundary conditions are imposed on Σ , similar to those analyzed previously by Dain (2004), by Maxwell (2004, 2005), and by Holst and Tsogtgerel (2013) as a model of black holes in various CMC settings, and by Holst, Meier, and Tsogtgerel (2013) in the setting of far-from-CMC solutions on compact manifolds with boundary. These "marginally trapped surface" Robin conditions ensure that the expansion scalars along null geodesics perpendicular to the boundary region Σ are non-positive, which is considered the correct mathematical model for black holes in the context of the Einstein constraint equations. Assuming a suitable form of weak cosmic censorship, the results presented in this article guarantee the existence of initial data that will evolve into a space-time containing an arbitrary number of black holes. A particularly important feature of our results are the minimal restrictions we place on the mean curvature, giving both near- and far-from-CMC results that are new.

CONTENTS

| 1. | Introduction | 1 |
|--|---|----|
| 2. | Asymptotically Euclidean Manifolds and Harmonic Functions | 4 |
| 3. | Overview of the Main Results | 6 |
| 4. | Properties of Linear Operators on Weighted Sobolev Spaces | 8 |
| 5. | Barriers for the Hamiltonian constraint | 9 |
| 6. | Non-CMC Solutions: Fixed Point Argument | 16 |
| 7. | Near-CMC Solutions: An Implicit Function Theorem Argument | 18 |
| Appendix A. Solutions to Semilinear, Boundary Value problems | | 20 |
| References | | 24 |

1. Introduction

In this paper we consider the Einstein constraint equations on an n-dimensional, asymptotically Euclidean manifold \mathcal{M} with boundary Σ . Using the recent work in [5, 11, 8], we show that far-from-CMC and near-CMC solutions exist to the conformal formulation of the Einstein constraints when nonlinear Robin boundary conditions are imposed on Σ similar to those developed in [3, 9, 8]. These "marginally trapped surface", Robin

Date: March 19, 2014.

Key words and phrases. Einstein constraint equations, weak solutions, asymptotically Euclidean, non-constant mean curvature, conformal method, manifolds with boundary.

MH was supported in part by NSF Awards 1065972, 1217175, and 1262982.

CM was supported in part by NSF Award 1065972.

conditions ensure that the expansion scalars along null geodesics perpendicular to the boundary region Σ are non-positive. Therefore, assuming a suitable form of weak cosmic censorship, the results presented here provide a method to construct initial data that will evolve into a space-time containing an arbitrary number of black holes. Moreover, this method imposes very few restrictions on the mean curvature.

We recall that the Einstein constraint equations on a given manifold $\mathcal M$ take the form

$$\hat{R} - \hat{K}^{ab}\hat{K}_{ab} + \hat{K} = \hat{\rho},\tag{1.1}$$

$$\hat{D}^{a}\hat{K} - \hat{D}_{b}\hat{K}^{ab} = -\hat{J}^{a},\tag{1.2}$$

where (1.1) is the Hamiltonian constraint and (1.2) is the momentum constraint. In the above system, \hat{R} and \hat{D} are the scalar curvature and connection with respect to the metric \hat{g}_{ab} , and \hat{K}_{ab} and \hat{K} are the extrinsic curvature tensor and its trace. The above underdetermined system imposes conditions on initial data $(\mathcal{M}, \hat{g}_{ab}, \hat{K}_{ab})$ for the initial value formulation of Einstein's equation.

In order to obtain solutions to (1.1)-(1.2) satisfying the marginally trapped surface conditions, we impose boundary conditions on $(\hat{g}_{ab}, \hat{K}_{ab})$ over Σ . Following the discussion in [9] and [8], a marginally trapped surface is one whose expansion along the incoming and outgoing orthogonal, null geodesics is non-positive. On the boundary Σ , the expansion scalars are given by

$$\hat{\theta}_{\pm} = \mp (n-1)\hat{H} + \operatorname{tr}_{\hat{g}}\hat{K} - \hat{K}(\hat{\nu}, \hat{\nu}), \tag{1.3}$$

where $(n-1)\hat{H} = \operatorname{div}_{\hat{g}}\hat{\nu}$ is the mean extrinsic curvature of Σ and $\hat{\nu}$ is the outward pointing, unit normal vector field to \mathcal{M} . Therefore, the surface Σ is called a marginally trapped surface if $\hat{\theta}_{\pm} \leq 0$. See [3, 9, 12] for details.

The problem we are interested in is to obtain solutions to the Einstein constraints for which $\theta_{\pm} \leqslant 0$. In order to formulate this problem as a determined system, we use the conformal method of Lichnerowicz, Choquet-Bruhat and York and the boundary conditions developed in [8]. Using the conformal method, one can transform (1.1)-(1.2) into a determined elliptic system by freely specifying conformal data, which consists of a Riemannian manifold (\mathcal{M},g) , a transverse traceless tensor σ , a mean curvature function τ , a non-negative energy density function ρ , and a vector field J. The Einstein constraints then become

$$-\Delta \phi + c_n R \phi + b_n \tau^2 \phi^{N-1} - c_n |\sigma + \mathcal{L}W|^2 \phi^{-N-1} - c_n \rho \phi^{-\frac{N}{2}} = 0, \tag{1.4}$$

$$\Delta_{\mathbb{L}}W + \frac{n-1}{n}\nabla\tau\phi^N + J = 0, \tag{1.5}$$

where ϕ is an undetermined positive scalar and W is an undetermined vector field. In the above equations, R is the scalar curvature of g, \mathcal{L} is the conformal Killing operator defined by

$$(\mathcal{L}W)_{ij} = D_i W_j + D_j W_i - \frac{2}{n} \nabla^k W_k g_{ij},$$

 ∇ and Δ are the connection and Laplacian associated with g, and $\Delta_{\mathbb{L}} = -\text{div} \circ \mathcal{L}$ is the vector Laplacian. The constants N, c_n and b_n are dimensional constants given by

$$N = \frac{2n}{n-2}$$
, $c_n = \frac{n-2}{4(n-1)}$, $b_n = \frac{n-2}{4n}$.

Combining the conformal method with the boundary conditions on θ_{\pm} in (1.3), one obtains the boundary conditions given in [8]. In particular, we will be interested in the

case when $\hat{\theta}_- = \theta_- \leqslant 0$ is freely specified. In this case, the boundary conditions in [8] are

$$\partial_{\nu}\phi + d_{n}H\phi + \left(d_{n}\tau - \frac{d_{n}}{n-1}\theta_{-}\right)\phi^{\frac{N}{2}} - \frac{d_{n}}{n-1}S(\nu,\nu)\phi^{-\frac{N}{2}} = 0 \quad \text{on } \Sigma, \tag{1.6}$$

$$(\mathcal{L}\mathbf{w})(\nu,\cdot) = \mathbf{V} \quad \text{on } \Sigma, \tag{1.7}$$

$$S(\nu,\nu) = \mathbf{V}(\nu) + \sigma(\nu,\nu) = ((n-1)\tau + |\theta_-|)\psi^N \geqslant 0 \quad \text{on } \Sigma.$$
 (1.8)

In (1.6), H is the rescaled extrinsic curvature for the boundary, $\nu = \phi^{\frac{N}{2}-1}\hat{\nu}$ is the rescaled normal vector field, and $d_n = \frac{n-2}{2}$ is a dimension dependent constant. The operators ∂_{ν} and \mathcal{L} are defined with respect to the specified metric g. In order to guarantee that $\theta_+ \leqslant 0$, the scalar function ψ is chosen so that $\phi \leqslant \psi$. In general, we are interested in solving the coupled conformal system (1.4)-(1.5) with the boundary conditions (1.6)-(1.8). We will refer to the boundary conditions (1.6)-(1.8) with the added condition that $\phi \leqslant \psi$ on $\mathcal M$ as marginally trapped surface boundary conditions, or more simply as marginally trapped surface conditions.

Our problem can now be expressed as a nonlinear, elliptic system of equations with Robin boundary conditions that is of the form

$$-\Delta \phi + c_n R \phi + b_n \tau^2 \phi^{N-1} - c_n |\sigma + \mathcal{L}W|^2 \phi^{-N-1} - c_n \rho \phi^{-\frac{N}{2}} = 0 \text{ on } \mathcal{M},$$

$$\partial_{\nu} \phi + d_n H \phi + \left(d_n \tau - \frac{d_n}{n-1} \theta_- \right) \phi^{\frac{N}{2}} - \frac{d_n}{n-1} S(\nu, \nu) \phi^{-\frac{N}{2}} = 0 \text{ on } \Sigma,$$

$$\Delta_{\mathbb{L}}W + \frac{n-1}{n} \nabla \tau \phi^N + J = \mathbf{0} \text{ on } \mathcal{M},$$

$$(\mathcal{L}\mathbf{w})(\nu, \cdot) = \mathbf{V} \text{ on } \Sigma.$$

$$(1.9)$$

One solves (1.9) for (ϕ, W) and then constructs a solution to the constraints from

$$\hat{g}_{ab} = \phi^{\frac{4}{n-2}} g_{ab}, \quad \hat{\rho} = \phi^{-\frac{3}{2}N+1} \rho, \quad \hat{J} = \phi^{-N} J,$$

$$\hat{K}^{ab} = \phi^{2\frac{(n+2)}{(n-2)}} (\sigma + \mathcal{L}W)^{ab} + \frac{\tau}{n} \phi^{\frac{4}{n-2}} g^{ab}.$$
(1.10)

If $\phi \leqslant \psi$, the expansion scalars associated with (\hat{g}, \hat{K}) will satisfy $\theta_{\pm} \leqslant 0$. In this case, (\hat{g}, \hat{K}) will be a solution to the coupled system which satisfies the marginally trapped surface conditions.

Boundary value problems similar to (1.9) were first studied in the constant mean curvature or CMC case. In [3] and [9], Dain and Maxwell proved the existence of apparent horizon solutions in this setting, with slight variations on the boundary condition (1.6). Then in [8], Holst and Tsogtgerel assembled a general collection of boundary conditions leading to marginally trapped surfaces that included the conditions of Maxwell and Dain, and then proved the existence of solutions to the Lichnerowicz problem on compact manifolds with boundary with simplifications of these conditions. It is important to note that the conditions in [8] imply an additional coupling between W and ϕ on the boundary, so even in the constant mean curvature case, the equations do not decouple. Holst and Tsogtgerel intentionally ignored this coupling in order to develop results for the Lichnerowicz equation alone as the first step in a program for the coupled system, and therefore did not construct solutions to the constraints satisfying the marginally trapped surface conditions. Their work then provided the mathematical framework for [11], where Holst, Meier, and Tsogtgerel showed that non-CMC solutions to the constraints exist satisfying the marginally trapped surface boundary conditions.

Outline of the Paper. The remainder of the paper is organized as follows. In Section 2, we introduce some basic notation and terminology in order to allow us to give a fairly complete overview of the main results in Section 3. We then develop some further notation and some basic supporting results in Section 4. The critical barrier (sub- and supersolution) constructions needed for our main results are then given in Section 5. The Schauder-based fixed-point framework is outlined in Section 6, followed by a proof of our main far-from-CMC result. A separate near-CMC result is then given in Section 7, based on the Implicit Function Theorem rather than a fixed-point argument. Some supporting results we need that supplement existing literature on this problem are given in Appendix A.

2. ASYMPTOTICALLY EUCLIDEAN MANIFOLDS AND HARMONIC FUNCTIONS

In this section, we introduce some basic notation and terminology in order to give an overview of the main results in Section 3. We will develop some further notation and some basic supporting results in Section 4 before giving the proofs of the main results in Sections 5–7.

Asymptotically Euclidean Manifolds. An n-dimensional, asymptotically Euclidean manifold (\mathcal{M}, q) is a non-compact Riemannian manifold, possibly containing a boundary, that can be decomposed into a compact set K and a finite number of ends E_1, \dots, E_k . Each E_i is diffeomorphic to the exterior of a ball in \mathbb{R}^n , and on each end the metric g tends towards the Euclidean metric g_E .

To formalize this definition, we recall the definition of the weighted Sobolev space $W^{k,p}_{\delta}(\mathcal{M})$ of scalar functions. (See [1] for an in depth discussion.) For $k \in \mathbb{N}$, $p \ge 1$, a given function $u \in W^{k,p}_{\delta}$ if

$$||u||_{W^{k,p}_{\delta}} = \sum_{|\beta| \leqslant k} ||r^{\delta - \frac{n}{p} + |\beta|} \partial^{\beta} u||_{L^p} < \infty.$$

$$(2.1)$$

In the above norm, partial derivatives are taken with respective to a fixed coordinate chart and r is a smooth positive function that agrees with |x| on each end E_i . For example, we may take $r(x) = \sqrt{1 + D(x, p_0)^2}$, where $D(x, p_0)$ denotes the distance from x to an arbitrary fixed point $p_0 \in K$. We will also consider the space of weighted, continuous functions $C^k_{\delta}(\mathcal{M})$, whose norm is given by

$$||u||_{C^k_\delta} = \sum_{|\alpha| \le k} \sup_{x \in \mathcal{M}} (r^{-\delta + |\alpha|} |\partial^{\alpha} u|).$$

The weighted Sobolev spaces and continuous spaces are related by the continuous embedding $W^{k,p}_{\delta} \hookrightarrow C^0_{\delta}$, which holds if k > n/p.

If $T^{a_1 a_2, \cdots a_r}_{b_1 b_2, \cdots b_s}$ is an (r, s)-tensor, we may define the point value of T by

$$|T| = (T_{b_1b_2,\cdots b_s}^{a_1a_2,\cdots a_r} T_{a_1a_2,\cdots a_r}^{b_1b_2,\cdots b_s})^{\frac{1}{2}}.$$

The above norms can then be applied to |T|, which allows one to consider weighted spaces $W^{k,p}_\delta(T^r_s\mathcal{M})$ of (r,s)-tensors. In particular, we let $\mathbf{W}^{k,p}_\delta=W^{k,p}_\delta(T\mathcal{M})$ denote the weighted Sobolev space of vector fields on \mathcal{M} .

We say that g tends towards g_E and is $W_{\delta}^{k,p}$ -asymptotically Euclidean if, for some $\delta < 0$.

$$g - g_E \in W^{k,p}_{\delta}. \tag{2.2}$$

We note that if $g-A_i^{N-2}g_E\in W^{k,p}_\delta(E_i)$ for each E_i , g is also $W^{k,p}_\delta$ -asymptotically Euclidean given that it will satisfy (2.2) with an appropriate change of coordinates (cf

[5]). Using these weighted spaces, we define an asymptotically Euclidean data set. As in [5], we say the data set $(\mathcal{M}, g, K, \rho, J)$ is asymptotically Euclidean if for some $\delta < 0$, $g - g_E \in W^{k,p}_{\delta}, K \in W^{k-1,p}_{\delta-1}$, and $\rho, J \in W^{k-2,p}_{\delta-2}$.

In the event that (\mathcal{M}, g) has a boundary Σ , we consider the Sobolev spaces $W^{k,p}(\Sigma)$ for $k \in \mathbb{N}$ and p > 1. These Banach spaces consist of the set of all functions u such that

$$||u||_{k,p;\Sigma} = \sum_{l \le k} ||\nabla^l u||_{p;\Sigma} < \infty,$$
 (2.3)

where the connection ∇ and integration are with respect to the boundary metric induced by q. This definition can be extended to obtain the fractional order Sobolev spaces $W^{s,p}(\Sigma)$ with $s \in \mathbb{R}$. See [7] for more details, including general results concerning multiplication properties of these spaces.

Asymptotic Limits and Harmonic Functions. We will seek solutions (ϕ, W) to the conformal equations where ϕ has fairly general asymptotic behavior. The following framework for representing this behavior is a generalization of the approach developed in [5], suitable for our needs here.

Given constants A_1, \dots, A_k , we seek solutions such that $\phi \to A_i$ on each end E_i . Let \mathcal{H} denote the space of smooth, harmonic functions with zero Neumann boundary conditions on Σ . By Proposition A.3, there exists a unique $\omega \in \mathcal{H}$ such that $\omega \to A_i$ on E_i . Therefore $\mathcal{H} \cong \mathbb{R}^k$ and if $\gamma < 0$, $\phi - \omega \in W^{2,p}_{\gamma}$ implies that $\phi \to A_i$ on each end. So ω encodes the asymptotic behavior of ϕ .

Because we can represent the asymptotic behavior of our solution ϕ by an element in $\omega \in \mathcal{H}$, we will seek solutions of the form $\phi = \omega + u$, where $u \in W^{2,p}_{\gamma}$. Therefore we define the space

$$\mathcal{H} + W_{\delta}^{k,p} = \{ \omega + u \mid \omega \in \mathcal{H}, \ u \in W_{\delta}^{k,p} \}.$$

If C^0 denotes the space of continuous functions on \mathcal{M} , we note that if k > n/p there exists a compact embedding

$$\mathcal{H} + W_{\delta}^{k,p} \hookrightarrow C^0, \tag{2.4}$$

given that $\mathbb{R}^k \oplus W^{k,p}_\delta \hookrightarrow \mathbb{R}^k \oplus C^0$ compactly and $\mathcal{H} + C^0 \subset C^0$. We will also need a way to compare the asymptotic limits of two functions f,g. We say f is **asymptotically bounded below** by g if

$$\lim_{|x|\to\infty} f \geqslant \lim_{|x|\to\infty} g \quad \text{ on each } E_i,$$

and f is **asymptotically bounded above** by g if g is asymptotically bounded below by f. Finally, given $g \leq h$ we say that f is **asymptotically bounded** by g and h if f is asymptotically bounded below by q and asymptotically bounded above by h.

Yamabe Invariant on Asymptotically Euclidean Manifolds with Boundary. To finish the discussion of notation needed for stating our main results, let us recall the definition of the Yamabe invariant on asymptotically Euclidean manifolds \mathcal{M} with boundary Σ . Define the following functional for compactly supported functions $f \in C_c^{\infty}$:

$$Q_g(f) = \frac{\int_{\mathcal{M}} |\nabla f|^2 + c_n R f^2 \, dV + \int_{\Sigma} d_n H f^2 \, dA}{\|f\|_{L^{\frac{2n}{n-2}}}^2}.$$
 (2.5)

Then as in [9], the Yamabe invariant on \mathcal{M} is

$$\mathcal{Y}_g = \inf_{f \in C_\infty^\infty(\mathcal{M}), \ f \neq 0} Q_g(f). \tag{2.6}$$

3. Overview of the Main Results

The main results for this paper concern the existence of far-from-CMC and near-CMC solutions to the conformal formulation of the Einstein constraint equations on an asymptotically Euclidean, n-dimensional manifold \mathcal{M} with compact boundary Σ . We assume that the boundary consists of m distinct components

$$\Sigma = \bigcup_{1}^{m} \Sigma_{i}, \quad \Sigma_{i} \cap \Sigma_{j} = \emptyset. \tag{3.1}$$

Here, each component Σ_i represents a marginally trapped surface, and \mathcal{M} is an embedded submanifold of some manifold \mathcal{N} . We view \mathcal{M} as the result of excising trapped regions C_i with boundary Σ_i from \mathcal{N} . Therefore, the following theorems provide conditions under which we may obtain solutions to the Einstein constraints outside of the singular trapped regions C_i with minimal assumptions on the mean curvature τ .

Our first Theorem is a far-from-CMC result in that it places no restrictions on the mean curvature function τ . However, to compensate for this assumption we require smallness assumptions on the other data.

Theorem 3.1. (Far-From-CMC) Suppose that (\mathcal{M}, g) is asymptotically Euclidean of class $W_{\gamma}^{2,p}$ with p > n and $2 - n < \gamma < 0$. Assume that $2 - n < \delta < \gamma/2$, and the data satisfies:

- $g \in \mathcal{Y}^+$, $\tau \in W^{1,p}_{\delta-1}$, $\sigma \in W^{1,2p}_{\delta-1}$ with $\|\sigma\|_{L^{\infty}_{\delta-1}}$ sufficiently small,
- ullet $ho \in L^{\infty}_{\gamma-2}$ with $\|\rho\|_{L^{\infty}_{\delta-2}}$ sufficiently small,
- $J \in \mathbf{L}_{\delta-2}^{p}$ with $\|J\|_{L_{\delta-2}^{p}}^{p-2}$ sufficiently small,
- $\theta_{-} \in W^{1-\frac{1}{p},p}(\Sigma), \, \theta_{-} < 0.$
- $\mathbf{V} \in \mathbf{W}^{1,p}$, $\mathbf{V}|_{\Sigma} = \left(((n-1)\tau + |\theta_{-}|/2)\psi^{N} \sigma(\nu, \nu) \right) \nu$, $((n-1)\tau + |\theta_{-}|/2) > 0$ and $||(n-1)\tau + |\theta_{-}|/2||_{W^{1-\frac{1}{p},p}(\Sigma)}$ sufficiently small.

Then on each end E_i there exists an interval $\mathcal{I}_i \subset (0, \infty)$ such that if $A_i \in \mathcal{I}_i$ are freely specified constants and ω is the associated harmonic function, there exists a solution (ϕ,W) to the conformal equations with boundary conditions (1.6)-(1.8) such that $\phi-\omega\in$ $W_{\gamma}^{2,p}$ and $W \in W_{\delta}^{2,p}$. Moreover, the function ψ can be chosen so that (ϕ, W) satisfies the marginally trapped surface boundary conditions.

Proof. The proof is given in Section 6.

The following Theorem complements Theorem 3.1 by showing that smallness assumptions on τ replace the need for smallness assumptions on σ and ρ . Given that the proof relies on the Implicit Function Theorem, solutions will be unique in this case.

Theorem 3.2. (Near-CMC with $g \in \mathcal{Y}^+$) Suppose that (\mathcal{M}, g) is asymptotically Euclidean of class $W_{\gamma}^{2,p}$ with p > n and $2 - n < \gamma < 0$. Assume that $2 - n < \delta < \gamma/2$, and the data satisfies:

- $q \in \mathcal{Y}^+$,
- $\|\tau\|_{W^{1,p}_{\delta-1}}$ is sufficiently small, and $\tau \geqslant 0$ on Σ ,
- $\sigma \in W_{\gamma-1}^{1,2p}$, $\rho \in L_{\gamma-2}^p$,
- $||J||_{\mathbf{L}^p_{\delta-2}}$ is sufficiently small,
- $\bullet \ \theta_-=0,$
- $\mathbf{V} \in \mathbf{W}^{1,p}$, $\mathbf{V}|_{\Sigma} = (((n-1)\tau)\phi^N \sigma(\nu,\nu))\nu$.

Then if $A_i \in (0,\infty)$ are freely specified constants on each end E_i and ω is the associated harmonic function, there exists a unique solution (ϕ, W) to the conformal equations with marginally trapped surface boundary conditions such that $\phi - \omega \in W^{2,p}_{\gamma}$ and $W \in W^{2,p}_{\delta}$.

Proof. The proof follows from Corollary 7.3 given in Section 7.

Our final Theorem states that we may replace the assumption that $q \in \mathcal{Y}^+$ with the assumption that R and H are bounded from below in terms of τ and $|\theta_-|$. In this case smallness assumptions are imposed on $\|\nabla \tau\|_{L^p_{\delta-1}}$ and $(2(n-1)\tau + |\theta_-|)$ on Σ .

Theorem 3.3. (Near-CMC with bounded R and H) Suppose that (\mathcal{M},g) is asymptotically Euclidean of class $W_{\gamma}^{2,p}$ with p>n and $2-n<\gamma<0$. Assume that $2-n < \delta < \gamma/2$, and the data satisfies:

- $\|\nabla \tau\|_{L^p_{\delta-2}}$ is sufficiently small, $\sigma \in W^{1,2p}_{\gamma-1}$, $\rho \in L^p_{\gamma-2}$, $J \in \mathbf{L}^p_{\delta-2}$,

- $\theta_{-} \in W^{1-\frac{1}{p},p}(\Sigma), \theta_{-} < 0.$
- $\mathbf{V} \in \mathbf{W}^{1,p}$, $\mathbf{V}|_{\Sigma} = (((n-1)\tau + |\theta_-|/2)\psi^N \sigma(\nu,\nu))\nu$,
- $(2(n-1)\tau + |\theta_-|) > 0$ is sufficiently small on Σ .

Let $A_i \in [1, \infty)$ be freely specified constants on each end E_i and let ω be the associated harmonic function. Then if

- $-c_n R \leqslant b_n \tau^2$ on $\{x \in \mathcal{M} : R(x) < 0\}$,
- $-H \le (\tau + |\theta_-|/(n-1))$ on $\{x \in \Sigma : H(x) < 0\}$,

there exists a solution (ϕ, W) to the conformal equations with boundary conditions (1.6)-(1.8) such that $\phi - \omega \in W^{2,p}_{\gamma}$ and $W \in W^{2,p}_{\delta}$. Moreover, the function ψ can be chosen so that (ϕ, W) satisfies the marginally trapped surface boundary conditions.

Proof. The proof is given in Section 6.

Remark 3.4. The conditions that $-c_nR \leq b_n\tau^2$ on $\{x \in \mathcal{M} : R(x) < 0\}$ and $\|\tau+|\theta_-|/(n-1)\|_{1-\frac{1}{p},p;\Sigma}$ and $\|\nabla\tau\|_{L^p_{\delta-1}}$ be sufficiently small place restrictions on the metric g. Namely, this method might not be applicable for metrics g which have large, negative scalar curvature. Similarly, the condition that $-H \leq (\tau + |\theta_-|/(n-1))$ on $\{x \in \Sigma : H(x) < 0\}$ and $\tau + |\theta_-|/(n-1)$ be small on Σ imposes conditions on the boundary. It is possible that these boundedness conditions on R and H relate to the positive Yamabe condition, however this relationship is not well understood (cf. [2]).

Remark 3.5. In Theorems 3.1-3.3 we assume the existence of a vector field $\mathbf{V} \in \mathbf{W}^{1,p}$ which satisfies

$$V|_{\Sigma} = \left(((n-1)\tau + |\theta_-|/2)\psi^N - \sigma(\nu, \nu) \right) \nu.$$

In the proof of Proposition 7.1 we explicitly construct a vector field satisfying these assumptions.

Theorem 3.2 follows from a variation of the Implicit Function Theorem argument developed in [2], where one perturbs τ and J from zero to obtain a small neighborhood of solutions about a known solution to the decoupled conformal equations. The proofs of Theorems 3.1 and 3.3 follow from a variation of the Schauder fixed point argument developed in [6, 7] for compact manifolds. This approach was adapted to asymptotically Euclidean manifolds in [5], and we use a variation of that argument. We briefly outline this method of proof below.

For a fixed W, we let $\mathcal{N}(\phi, W)$ denote the Lichnerowicz operator on the left of (1.4) with boundary operator on the left of (1.6). In this notation, a solution to the coupled system (1.9) satisfies $\mathcal{N}(\phi, \mathcal{S}(\phi)) = 0$, where $\mathcal{S}(\phi)$ denotes the solution to the momentum constraint (1.5) with boundary conditions (1.7) for a given ϕ .

A rough outline of the Schauder fixed point argument is as follows. If C^0_+ denotes the spaces of positive, continuous functions, for a given $\phi \in \mathcal{C}^0_+$ and $\psi \in W^{2,p}(\Sigma)$ let $W = \mathcal{S}(\phi)$ denote the momentum constraint solution map with boundary condition (1.7)-(1.8). Similarly, for a given $W \in W^{2,p}_\delta$ and sub-and supersolutions $\phi_- \leqslant \phi_+$, Theorems A.4 and A.5 in Appendix A imply that for a given $\omega \in \mathcal{H}$ which is asymptotically bounded by ϕ_- and ϕ_+ , there exists a unique solution to $\mathcal{N}(\phi,W)=0$ such that $\phi-\omega\in W^{2,p}_\gamma$. Therefore we let $\phi=T(W)$ denote the Hamiltonian constraint solution map with boundary conditions (1.6). If i is the compact embedding $\mathcal{H}\oplus W^{2,p}_\gamma\hookrightarrow C^0$ defined in (2.4), then we set

$$\mathcal{G}(\phi) = i(T(\mathcal{S}(\phi))). \tag{3.2}$$

A solution to the coupled system with the specified boundary conditions will be a fixed point of this map. In order to apply the Schauder fixed point argument in [7], we must show that this map is compact and invariant on a certain subset of C^0_+ .

The primary difficulty in applying this fixed point argument is in constructing the closed, bounded, and convex subset of C^0 on which the map $\mathcal{G}(\phi)$ is invariant. The construction of this set requires global sub- and supersolutions ϕ_- and ϕ_+ of the Hamiltonian constraint, and once these are obtained the process is fairly straightforward. See [7, 11, 5]. In Section 5, we will construct global sub-and super-solutions for the Hamiltonian constraint and in Section 6 we use this framework to prove Theorems 3.1 and 3.3. Then in Section 7 we use the Implicit Function Theorem to obtain the near-CMC results in Theorem 3.2.

4. Properties of Linear Operators on Weighted Sobolev Spaces

Using the definition of Yamabe invariant given in 2.5, we compile some useful facts from [9] about the operators

$$\mathcal{P}_1 = (-\Delta + c_n R, \partial_{\nu} + d_n H)$$
 and $\mathcal{P}_2 = (-\Delta_{\mathbb{L}}, B),$

where $BW=\mathcal{L}W(\nu,\cdot)$. In the following proposition, we summarize the properties of both \mathcal{P}_1 and \mathcal{P}_2 . We write $L^p_{\delta-2}\times W^{1-\frac{1}{p},p}$ to indicate $L^p_{\delta-2}(\mathcal{M})\times W^{1-\frac{1}{p},p}(\Sigma)$ in the case of \mathcal{P}_1 and $L^p_{\delta-2}(T\mathcal{M})\times W^{1-\frac{1}{p},p}(T\Sigma)$ in the case of \mathcal{P}_2 .

Proposition 4.1. Suppose (\mathcal{M},g) is asymptotically Euclidean of class $W_{\gamma}^{k,p}$ with $k\geqslant 2$, k>n/p, and $2-n<\delta<0$. Then $\mathcal{P}_i:W_{\delta}^{2,p}\to L_{\delta-2}^p\times W^{1-\frac{1}{p},p}$ is Fredholm with index zero. Moreover, if $\mathcal{Y}_g>0$, then \mathcal{P}_1 is an isomorphism and if p>n or \mathcal{M} possesses no conformal Killing fields, then \mathcal{P}_2 is an isomorphism. Finally, if \mathcal{P}_i is an isomorphism and $\mathcal{P}_iv=(f,g)\in L_{\delta-2}^p\times W^{1-\frac{1}{p},p}$, then there exists C>0 such that the following estimate is satisfied:

$$||v||_{W^{2,p}_{\delta}} \leqslant C \left(||f||_{L^{p}_{\delta-2}} + ||g||_{W^{1-\frac{1}{p},p}} \right).$$
 (4.1)

Proof. See Proposition 1, Proposition 3, Proposition 6 and Theorem 3 in [9]. \Box

With Proposition 4.1 in hand, we can now prove the following important estimate in the case when k = 2. This result is based on a similar estimate in [4].

Proposition 4.2. Suppose (\mathcal{M},g) is asymptotically Euclidean of class $W^{2,p}_{\gamma}$ with n < p, and let r be the function defined in (2.1). Then for a given $\phi \in L^{\infty}_+$, if $W \in W^{2,p}_{\delta}$ is solves the momentum constraint (1.5) with boundary conditions (1.7), where $2 - n \leq \delta < 0$, there exists C > 0 such that the following estimate holds:

$$\|\mathcal{L}W\|_{\infty} \leqslant Cr^{\delta-1} \left(\|\nabla \tau\|_{L^{p}_{\delta-2}} \|\phi\|_{\infty}^{N} + \|J\|_{L^{p}_{\delta-2}} + \|\mathbf{V}\|_{W^{1-\frac{1}{p},p}(T\Sigma)} \right)$$
(4.2)

Proof. By Proposition 4.1 we have

$$||W||_{W^{2,p}_{\delta}} \le c \left(||\nabla \tau \phi^{N}||_{L^{p}_{\delta-2}} + ||J||_{L^{p}_{\delta-2}} + ||\mathbf{V}||_{W^{1-\frac{1}{p},p}(T\Sigma)} \right). \tag{4.3}$$

The continuous embedding $W^{1,p}_{\delta-1} \hookrightarrow C^0_{\delta-1}$ implies that

$$\|\mathcal{L}W\|_{C^0_{\delta-1}} \leqslant C_1 \|LW\|_{W^{1,p}_{\delta-1}} \leqslant C_2 \|W\|_{W^{2,p}_{\delta}},$$

and combining this with estimate (4.3) we have

$$\|\mathcal{L}W\|_{C^{0}_{\delta-1}} \leqslant C\left(\|\phi\|_{\infty}^{N} \|\nabla \tau\|_{L^{q}_{\delta-2}} + \|J\|_{L^{p}_{\delta-2}} + \|\mathbf{V}\|_{W^{1-\frac{1}{p},p}(T\Sigma)}\right). \tag{4.4}$$

The above estimate and the definition of the $C_{\delta-1}^0$ norm imply the result.

Propositions 4.1 and 4.2 will be essential in determining our global barriers. In particular, Proposition 4.2 is our primary tool to control the point-wise values of the solution of the momentum constraint $W = \mathcal{S}(\phi)$ in terms of ϕ . This will be vital when we construct our global supersolution in the next section.

5. Barriers for the Hamiltonian constraint

A critical component of fixed-point arguments for nonlinear elliptic equations are the development of *a priori* estimates, and/or sub- and supersolutions. These so-called barriers are an essential component for building the \mathcal{G} -invariant set necessary for our fixed point argument, where we recall that \mathcal{G} is the nonlinear fixed point operator defined in (3.2). Therefore, in this section we will develop several global sub-and supersolution constructions for the Hamiltonian constraint equation (1.4).

If γ is the trace operator associated with Σ , we define the operators

$$A_{L}(\phi) = \begin{pmatrix} -\Delta\phi + c_{n}a_{R}\phi \\ \gamma(\partial_{\nu}\phi) + d_{n}H(\gamma\phi) \end{pmatrix},$$

$$F(\phi, W) = \begin{pmatrix} b_{n}\tau^{2}\phi^{N-1} - c_{n}|\sigma + \mathcal{L}W|^{2}\phi^{-N-1} - c_{n}\rho\phi^{-\frac{N}{2}} \\ (d_{n}\gamma\tau - \frac{d_{n}}{n-1}\theta_{-})(\gamma(\phi))^{\frac{N}{2}} - \frac{d_{n}}{n-1}S(\nu, \nu)(\gamma(\phi))^{-\frac{N}{2}} \end{pmatrix}.$$

The Hamiltonian constraint with boundary conditions (1.6)-(1.8) can be written succinctly as

$$\mathcal{N}(\phi, W) = A_L(\phi) + F(\phi, W) = 0. \tag{5.1}$$

Using this notation, we recall that for a given vector field W, if the functions ϕ_- and ϕ_+ satisfy

$$\mathcal{N}(\phi_-, W) \leqslant 0$$
 and $\mathcal{N}(\phi_+, W) \geqslant 0$,

then ϕ_{-} is called a subsolution and ϕ_{+} is a supersolution.

As in [7], to obtain a fixed point of the coupled conformal equations we require a slightly more restrictive class of sub- and supersolutions. If $W=W(\phi)$ denotes the solution of the momentum constraint for a given ϕ and

$$\mathcal{N}(\phi_-, W(\phi)) \leqslant 0$$
 for all $\phi \geqslant \phi_-$,

then ϕ_{-} is a global subsolution. Similarly, ϕ_{+} is a global supersolution if

$$\mathcal{N}(\phi_+, W(\phi)) \geqslant 0$$
 for all $\phi \leqslant \phi_+$.

In the following discussion we will require that when the vector field $W \in \mathbf{W}_{\delta}^{2,p}$ is given by the solution of the momentum constraint equation (1.5) with the source term $\phi \in L^{\infty}$,

$$|\mathcal{L}W|^2 \leqslant r^{2\delta - 2} (\mathbf{k}_1 \|\phi\|_{\infty}^{2N} + \mathbf{k}_2),$$
 (5.2)

with some positive constants k₁ and k₂. The following proposition justifies this bound.

Proposition 5.1. Let the assumptions of Proposition 4.2 hold with $p \in (n, \frac{2\alpha+1}{2})$, $\alpha > n$, and $s \in (1+\frac{(n-1)}{p}-\frac{(n-1)}{\alpha},1+\frac{(n-1)}{p})$. Suppose that W satisfies the momentum constraint (1.5) with boundary conditions (1.7)-(1.8), where $\mathbf{V}(\nu) = ((n-1)\tau + |\theta_-|/2)\psi^N$. Then W satisfies the bound (5.2) with

$$\begin{aligned} \mathbf{k}_{1} &= 2C_{1}^{2} \|\nabla \tau\|_{L_{\delta-2}^{p}}^{2}, \\ \mathbf{k}_{2} &= 2C_{2}^{2} \left(\|J\|_{L_{\delta-2}^{p}} + \|\sigma(\nu, \nu)\|_{1-\frac{1}{p}, p; \Sigma} \right. \\ &+ \left. \|((n-1)\tau + |\theta_{-}|/2)\|_{1-\frac{1}{p}, p; \Sigma} \|\psi\|_{\infty}^{N-1} \|\psi\|_{s, p; \Sigma} \right)^{2}. \end{aligned}$$

$$(5.3)$$

Proof. Proposition 4.2 implies that

$$\|\mathcal{L}W\|_{\infty}^{2} \leqslant C^{2} r^{2\delta-2} \left(\|\phi\|_{\infty}^{N} \|\nabla \tau\|_{L_{\delta-2}^{p}} + \|J\|_{L_{\delta-2}^{p}} + \|\mathbf{V}\|_{1-\frac{1}{p},p;\Sigma} \right)^{2}.$$
 (5.4)

If $\mathbf{V}(\nu) = ((n-1)\tau + |\theta_-|/2)\psi^N - \sigma(\nu,\nu)$ as in Theorem 3.1, then $\mathbf{V}|_{\Sigma} = \mathbf{X} + \mathbf{Y}$, where $\mathbf{X}(\nu) = \mathbf{V}(\nu)$ and $\mathbf{Y}(\nu) = \mathbf{0}$. In practice, we will assume that $\mathbf{Y} = \mathbf{0}$ (cf. Proposition 7.1), and we have that

$$\|\mathbf{V}\|_{1-\frac{1}{p},p;\Sigma} \leqslant C\left(\|((n-1)\tau + |\theta_{-}|/2)\psi^{N}\|_{1-\frac{1}{p},p;\Sigma} + \|\sigma(\nu,\nu)\|_{1-\frac{1}{p},p;\Sigma}\right). \tag{5.5}$$

We now apply Lemma A.21 from [7] with $\sigma=1-\frac{1}{p},\ p=q,$ and $s\in(1+\frac{(n-1)}{p}-\frac{(n-1)}{\alpha},1+\frac{(n-1)}{p}).$ This gives us that

$$\|((n-1)\tau + |\theta_{-}|/2)\psi^{N}\|_{1-\frac{1}{p},p;\Sigma}$$

$$\leq C \|((n-1)\tau + |\theta_{-}|/2)\|_{1-\frac{1}{p},p;\Sigma} (\|\psi^{N}\|_{\infty} + N\|\psi^{N-1}\|_{\infty}\|\psi\|_{s,p;\Sigma}) .$$

The embedding $W^{s,p}(\Sigma) \hookrightarrow L^{\infty}(\Sigma)$ implies that

$$\|((n-1)\tau+|\theta_-|/2)\psi^N\|_{1-\frac{1}{p},p;\Sigma}\leqslant C\|((n-1)\tau+|\theta_-|/2)\|_{1-\frac{1}{p},p;\Sigma}\|\psi\|_{\infty}^{N-1}\|\psi\|_{s,p;\Sigma},$$

which combined with (5.4) implies the result.

In the following discussion, we let $\omega \in \mathcal{H}$ be such that $\omega \to A_j > 0$ on each end E_j . Additionally, given any scalar function $u \in L^{\infty}$, we use the notation

$$u^{\wedge} := \operatorname{ess\ sup} u, \qquad u^{\vee} := \operatorname{ess\ inf} u.$$

We are now ready to construct our global barriers. The following Theorem provides conditions under which we can construct a global supersolution for $\mathcal N$ given that W satisfies the boundary conditions (1.7)-(1.8). For this particular construction, if we want to freely specify $\psi \in W^{1-\frac{1}{p},p}(\Sigma)$ we are required to assume that $(2(n-1)\tau+|\theta_-|)$ is sufficiently small on Σ .

Theorem 5.2. (Far-From-CMC Global Supersolution) Suppose that (\mathcal{M}, q) is asymptotically Euclidean of class $W_{\gamma}^{2,p}$, with n < p and $\gamma \in (2-n,0)$, and that 2-n < p $\delta < \gamma/2$. Additionally assume that $\mathcal{Y}_g > 0$, $\tau \in W^{1,p}_{\delta-1}$, and that $\sigma \in L^{\infty}_{\delta-1} \cap W^{1,2p}_{\delta-1}$ $J\in L^p_{\delta-2}$ and $\rho\in L^\infty_{\delta-2}$ are sufficiently small. Also assume that for $0<\psi\in W^{2,p}_\delta$, $(2(n-1)\tau + |\theta_-|)\psi^N > 0$ is sufficiently small on Σ . Then there exists a global supersolution $\phi_+ > 0$ to the Hamiltonian constraint with boundary condition (1.6)-(1.8) such that $\phi_+ - \beta \omega \in W^{2,p}_{\gamma}$ for some $\beta > 0$ sufficiently small.

Proof. Let $\Lambda \in L^p_{\gamma-2}$ be a positive function that agrees with $r^{\gamma-2}$ outside of a compact set and let $\lambda \in W^{1-\frac{1}{p},p}$ be a positive function on Σ . Then by Proposition 4.1 and Proposition 3.2 in [5] there exists solution $u \in W^{2,p}_{\gamma}$ solving

$$-\Delta u + c_n R u = \Lambda - c_n \omega R,$$

$$\partial_{\nu} u + d_n H u = \lambda - d_n \omega H.$$
(5.6)

Let $\phi_+ = \beta(u + \omega)$, where $\beta > 0$ will be determined. By the maximum principles A.1 and A.2 we have that $\phi_+>0$. We recall from Proposition 5.1 that we may bound $a_{\mathcal{L}W}$ in terms of the source function ϕ . Using this bound and the fact that $a_W = c_n |\sigma + \mathcal{L}W|^2 \leqslant$ $2|\sigma|^2 + 2|\mathcal{L}W|^2$, we obtain the bound

$$a_W^{\wedge} \leqslant r^{2\delta - 2} (K_1 \|\phi\|_{\infty}^{2N} + K_2),$$
 (5.7)

where

$$K_1 = C_1 \mathbf{k}_1$$
 and $K_2 = 2r^{2-2\delta} (\sigma^2)^{\wedge} + C_2 \mathbf{k}_2$,

and k_1 and k_2 are the same constants in (5.3). We let $W(\phi)$ denote a solution to the momentum constraint for a given $\phi < \phi_+$ and define $k_3 = (\frac{\sup \phi_+}{\inf \phi_+})^{2N}$. Applying the Hamiltonian constraint (5.1) to ϕ_+ and using the fact that $S(\nu,\nu) = ((n-1)\tau + |\theta_-|/2)\psi^N$, we obtain

$$\mathcal{N}(\phi_{+}, W(\phi))
= \begin{pmatrix} -\Delta\phi_{+} + c_{n}R\phi_{+} + b_{n}\tau^{2}\phi_{+}^{N-1} - a_{W}\phi_{+}^{-N-1} - c_{n}\rho\phi_{+}^{-N/2} \\ \partial_{\nu}\phi_{+} + d_{n}H\phi_{+} + (d_{n}\tau - \frac{d_{n}}{n-1}\theta_{-})\phi_{+}^{\frac{N}{2}} - \frac{d_{n}}{n-1}(((n-1)\tau + |\theta_{-}|/2)\psi^{N})\phi_{+}^{-\frac{N}{2}} \end{pmatrix}
\geqslant \begin{pmatrix} \beta\Lambda + \beta^{N-1}c_{n}\tau^{2}(u+\omega)^{N-1} - r^{2\delta-2}(K_{1}(\phi^{\wedge})^{2N} + K_{2})\phi_{+}^{-N-1} - c_{n}\rho\phi_{+}^{-N/2} \\ \beta\lambda + (d_{n}\tau - \frac{d_{n}}{n-1}\theta_{-})\phi_{+}^{\frac{N}{2}} - \frac{d_{n}}{n-1}(((n-1)\tau + |\theta_{-}|/2)\psi^{N})\phi_{+}^{-\frac{N}{2}} \end{pmatrix}
\geqslant \begin{pmatrix} \beta\Lambda - r^{2\delta-2}K_{1}k_{3}\beta^{N-1}(u+\omega)^{N-1} - K_{2}r^{2\delta-2}\phi_{+}^{-N-1} - c_{n}\rho\phi_{+}^{-N/2} \\ \beta\lambda + (d_{n}\tau - \frac{d_{n}}{n-1}\theta_{-})\phi_{+}^{\frac{N}{2}} - \frac{d_{n}}{n-1}(((n-1)\tau + |\theta_{-}|/2)\psi^{N})\phi_{+}^{-\frac{N}{2}} \end{pmatrix}.$$

As in Theorem 4.1 in [5], the decay rate on Λ ensures that we can choose β sufficiently small so that

$$\frac{\beta\Lambda}{2} - r^{2\delta - 2} K_1 \mathbf{k}_3 \beta^{N-1} (u + \omega)^{N-1} > 0.$$

The smallness assumptions on σ , ρ , J on \mathcal{M} and the smallness assumptions on $(2(n-1)^{n})^{-1}$ $1)\tau + |\theta_-|\psi^N$ on Σ imply that we can ensure that the first equation in the above array is nonnegative. For this fixed β , we observe that

$$\beta \lambda + (d_n \tau - \frac{d_n}{n-1} \theta_-) \phi_+^{\frac{N}{2}} > 0.$$

Therefore the smallness assumption on $(2(n-1)\tau + |\theta_-|)\psi^N$ implies that the boundary equation can be made nonnegative as well. Therefore, $\phi_+ = \beta(u+\omega)$ will be a global super-solution of the Hamiltonian constraint with boundary condition (1.6)-(1.8) if $\beta > 0$ is sufficiently small and the conformal data satisfies the above assumptions.

In the following Theorem, we show that when $\psi = \phi_+ = \beta(u + \omega)$, where u satisfies (5.6), ϕ_+ will be a global super-solution to the Hamiltonian constraint with boundary conditions (1.6)-(1.8) provided that $\beta>0$ is chosen sufficiently small and our data σ,ρ,J is sufficiently small. The significance of this result is that the supersolution acts as an a priori upper bound for the fixed point solution ϕ , and therefore $\phi\leqslant\phi_+=\psi$. This implies that the resulting fixed point ϕ will satisfy the marginally trapped surface conditions, which is why we refer to the following supersolution construction as a marginally trapped surface supersolution.

Theorem 5.3. (Marginally Trapped Surface Supersolution for $g \in \mathcal{Y}^+$) Let the assumptions of Theorem 5.2 hold with the exception of the smallness assumption on $(2(n-1)\tau + |\theta_-|)\psi^N$, and let u satisfy equation (5.6). Then there exists a $\beta > 0$ such that if $\phi_+ = \beta(u + \omega)$ and $\psi = \beta(u + \omega)$ on Σ , ϕ_+ will be a global supersolution to the Hamiltonian constraint with boundary condition (1.6)-(1.8) that also imposes the marginally trapped surface condition.

Proof. As in the proof of Theorem 5.2, we let $a_W = c_n |\sigma + \mathcal{L}W|^2 \leqslant 2|\sigma|^2 + 2|\mathcal{L}W|^2$ and apply the estimate (5.7), where the constants k_1, k_2, K_1 and K_2 are the same as in the previous proof. We apply the Hamiltonian constraint (5.1) to ϕ_+ and use the fact that $S(\nu, \nu) = ((n-1)\tau + |\theta_-|/2)\phi_+^N$ to obtain

$$\mathcal{N}(\phi_{+}, W(\phi))
= \begin{pmatrix} -\Delta\phi_{+} + c_{n}R\phi_{+} + b_{n}\tau^{2}\phi_{+}^{N-1} - a_{W}\phi_{+}^{-N-1} - c_{n}\rho\phi_{+}^{-N/2} \\ \partial_{\nu}\phi_{+} + d_{n}H\phi_{+} + (d_{n}\tau - \frac{d_{n}}{n-1}\theta_{-})\phi_{+}^{\frac{N}{2}} - \frac{d_{n}}{n-1}((n-1)\tau + |\theta_{-}|/2)\phi_{+}^{\frac{N}{2}} \end{pmatrix}
\geqslant \begin{pmatrix} \beta\Lambda + \beta^{N-1}b_{n}\tau^{2}(u+\omega)^{N-1} - r^{2\delta-2}(K_{1}(\phi^{\wedge})^{2N} + K_{2})\phi_{+}^{-N-1} - c_{n}\rho\phi_{+}^{-N/2} \\ \beta\lambda + \left(\frac{d_{n}}{2(n-1)}|\theta_{-}|\right)\phi_{+}^{\frac{N}{2}} \end{pmatrix}.$$

We now observe that

$$\mathbf{k}_{2} \leqslant 4C^{2} \left(\|J\|_{L_{\delta-2}^{p}} + \|\sigma(\nu,\nu)\|_{1-\frac{1}{p},p;\Sigma} \right)^{2}
+ 4C \left(\|((n-1)\tau + |\theta_{-}|/2)\|_{1-\frac{1}{p},p;\Sigma} \|\psi\|_{\infty}^{N-1} \|\psi\|_{s,p;\Sigma} \right)^{2}
= C_{1}(\sigma,J) + \beta^{2N} C_{2}(\tau,\theta_{-},u),$$
(5.8)

and we have that

$$K_2 \leq 2r^{2-2\delta}(\sigma^2)^{\wedge} + C_1(\sigma, J) + \beta^{2N}C_2(\tau, \theta_-, u) = C_3(\sigma, J) + \beta^{2N}C_2(\tau, \theta_-, u).$$

If k_3 is as in the proof of Theorem 5.2 and if we choose $\beta > 0$ sufficiently small so that

$$\frac{\beta\Lambda}{2} - \left(\mathbf{k}_3 K_1(u+\omega)^{-N-1} + C_2(\tau,\theta_-,u) \right) \beta^{N-1} r^{2\delta-2} > 0 \text{ on } \mathcal{M},$$

then we can ensure ϕ_+ will be a super-solution by imposing smallness assumptions on σ , ρ and J as in the proof of Theorem 5.2. Given that the second equation is positive for any choice of $\beta > 0$, $\phi_+ = \beta(u + \omega)$ will be a global super-solution of the Hamiltonian constraint with boundary condition (1.6) for $\beta > 0$ sufficiently small and conformal data satisfying the assumptions of the Theorem.

The following Lemma provides us with a method of constructing a global supersolution in the event that g is not in the positive Yamabe class. However, we require that the scalar curvature R and boundary mean curvature H be bounded by functions of τ and θ_- on the sets where they are negative.

Lemma 5.4. Suppose that (\mathcal{M},g) is asymptotically Euclidean of class $W^{2,p}_{\gamma}$ with p>n and $\gamma\in(2-n,0)$. Assume that $2-n<\delta<\gamma/2$ and that $\tau\in W^{1,p}_{\delta-1}$, $\rho\in L^p_{\gamma-2}$, $\sigma\in W^{1,2p}_{\gamma-1}$, $J\in \mathbf{L}^p_{\delta-2}$, $\theta_-\in W^{1-\frac{1}{p},p}(\Sigma)$, and $(2(n-1)\tau+|\theta_-|)>0$ on Σ . Additionally assume that

- $-c_n R \leqslant b_n \tau^2$ on $\{x \in \mathcal{M} : R(x) < 0\}$,
- $-H \le (\tau + |\theta_-|/(n-1))$ on $\{x \in \Sigma : H(x) < 0\}$,

and that $A \in L^p_{\gamma-2}$ is nonnegative. Then for $\omega \in \mathcal{H}$ with asymptotic limits $A_i \in [1, \infty)$, there exists a solution ϕ_A to the equation

$$-\Delta \phi_{A} = c_{n} A \phi_{A}^{-N-1} + c_{n} \rho \phi_{A}^{-\frac{N}{2}},$$

$$\partial_{\nu} \phi_{A} = \frac{d_{n}}{n-1} A \phi_{A}^{-\frac{N}{2}},$$
(5.9)

such that $\phi_A - \omega \in W^{2,p}_{\gamma}$. If W_{ϕ} denotes the solution to the momentum constraint for $\phi \leqslant \phi_A$, $A \geqslant |\sigma + \mathcal{L}W_{\phi}|^2$ on \mathcal{M} , and $A \geqslant S(\nu, \nu) = (\sigma(\nu, \nu) + \mathcal{L}W_{\phi}(\nu, \nu))$ on Σ , then ϕ_A will be a global supersolution of the Hamiltonian constraint with boundary conditions (1.6)-(1.8).

Proof. We note that $\phi_- = 1$ is a subsolution to (5.9). We obtain a supersolution by letting $\phi_+ = \beta(u+1)$, where $\beta \gg 1$ is sufficiently large and u is the solution to

$$-\Delta u = c_n A + c_n \rho$$

$$\partial_{\nu} u = \frac{d_n}{n-1} A.$$
(5.10)

The maximum principle implies that $\phi_+ > 0$ and for $\beta \gg 1$, we have $\phi_- \leqslant \phi_+$. We may then apply Theorem A.4 to obtain a solution $\phi_A \geqslant 1$ which tends to freely specified $A_i \in [1,\beta]$ on each end E_i . As β can be arbitrarily large, the asymptotic limits A_i can be freely specified numbers in $[1,\infty)$.

Now we compute $\mathcal{N}(\phi_A, W)$ for $\phi \leqslant \phi_A$, where $W = W(\phi)$ depends on ϕ . We obtain

$$\mathcal{N}(\phi_{A}, W) = \begin{pmatrix}
-\Delta \phi_{A} + c_{n} R \phi_{A} + b_{n} \tau^{2} \phi_{A}^{N-1} - c_{n} | \sigma + \mathcal{L}W |^{2} \phi_{A}^{-N-1} - c_{n} \rho \phi_{A}^{-\frac{N}{2}} \\
\partial_{\nu} \phi_{A} + d_{n} H \phi_{A} + \left(d_{n} \tau + \frac{d_{n}}{n-1} | \theta_{-} | \right) \phi_{A}^{\frac{N}{2}} - \frac{d_{n}}{n-1} S(\nu, \nu) \phi_{A}^{-\frac{N}{2}} \end{pmatrix} \\
\geqslant \begin{pmatrix}
c_{n} R \phi_{A} + b_{n} \tau^{2} \phi_{A}^{N-1} + c_{n} (A - | \sigma + \mathcal{L}W |^{2}) \phi_{A}^{-N-1} \\
d_{n} H \phi_{A} + \left(d_{n} \tau + \frac{d_{n}}{n-1} | \theta_{-} | \right) \phi_{A}^{\frac{N}{2}} + \frac{d_{n}}{n-1} (A - S(\nu, \nu)) \phi_{A}^{-\frac{N}{2}} \end{pmatrix} \geqslant 0,$$

where the last is expression is nonnegative given the assumptions on A, the fact that $\phi_A \geqslant 1$, and the assumption that $-c_n R \leqslant b_n \tau^2$ and $-H \leqslant (\tau + |\theta_-|/(n-1))$ on the sets where R and H are negative.

The previous Lemma tells us that ϕ_A will be a global supersolution to the Hamiltonian constraint with boundary conditions (1.6) provided that we can find a function A which bounds the terms in $\mathcal{N}(\phi, W(\phi))$ that depend on $W(\phi)$. We construct such a function in the following theorem and show that when $\psi = \phi_A$ the function ϕ_A is a marginally trapped surface supersolution.

Theorem 5.5. (Global Supersolution for bounded R and H) Let the assumptions of Lemma 5.4 hold and suppose that V satisfies the conditions of Theorem 3.3. Additionally assume that Σ is compact. Then there exists an $\epsilon > 0$ such that if

$$\|\nabla \tau\|_{L^p_{\delta-2}} < \epsilon \quad and \quad \|\tau + |\theta_-|\|_{W^{1-\frac{1}{p},p}(\Sigma)} < \epsilon, \tag{5.11}$$

and $A = Cr^{2\delta-2}$ for some constant C > 0, the solution ϕ_A to (5.9) will be a global supersolution to the Hamiltonian constraint with boundary conditions (1.6)-(1.8). Moreover, if $\phi_A = \psi$ the marginally trapped surface condition will hold.

Proof. Lemma 5.4 implies that $\phi_+ = \phi_A$ will be a supersolution for $\phi \leqslant \phi_A$ provided that we can choose $A \geqslant |\sigma + \mathcal{L}W|^2$ on \mathcal{M} and $A \geqslant (\sigma(\nu, \nu) + \mathcal{L}W(\nu, \nu))$ on Σ . Using the estimate from Proposition 5.1, we have

$$|\sigma + \mathcal{L}W|^2 \le 2|\sigma|^2 + 2|\mathcal{L}W|^2 \le 2|\sigma|^2 + 2r^{2\delta-2}(\mathbf{k}_1 ||\phi_A||_{\infty}^{2N} + \mathbf{k}_2),$$

where

$$\begin{split} \mathbf{k}_1 &= 2C_1^2 \|\nabla \tau\|_{L^p_{\delta-2}}^2, \\ \mathbf{k}_2 &= 2C_2^2 \left(\|J\|_{L^p_{\delta-2}} + \|\sigma(\nu,\nu)\|_{1-\frac{1}{p},p;\Sigma} \right. \\ &+ \left. \|((n-1)\tau + |\theta_-|/2)\|_{1-\frac{1}{p},p;\Sigma} \|\psi\|_{\infty}^{N-1} \|\psi\|_{s,p;\Sigma} \right)^2. \end{split}$$

Setting $A = Cr^{2\delta-2}$, if we can find C > 0 so that

$$2\|\sigma\|_{L_{\delta-1}^{\infty}}^{2} + 2(\mathbf{k}_{1}\|\phi_{A}\|_{\infty}^{2N} + \mathbf{k}_{2}) \leqslant C,$$

$$|S(\nu,\nu)| \leqslant (2(n-1)\tau + |\theta_{-}|)\psi^{N} \leqslant C \min_{r \in \Sigma} (r^{2\delta-2}),$$
(5.12)

the conditions of Lemma 5.4 will be satisfied and ϕ_A will be a global supersolution. For arbitrary an $\psi \in W^{1-\frac{1}{p},p}(\Sigma)$ that is independent of A, we choose

$$C>\max\{2(\|\sigma\|_{L^\infty_{\delta-1}}^2+\mathtt{k}_2),\alpha\max_{x\in\Sigma}(2((n-1)\tau+|\theta_-|)\psi^N)\},$$

where $\alpha=1/(\min_{x\in\Sigma}(r^{2\delta-2}))$. Taking $\mathbf{k}_1=\|\nabla\tau\|_{L^p_{\delta-2}}^2$ to be sufficiently small we can ensure that both inequalities in (5.12) hold.

To obtain our marginally trapped supersolution, we set $\psi = \phi_A$. In this case we take

$$C > 2(\|\sigma\|_{L^{\infty}_{\delta-1}}^2 + \mathbf{k}_2),$$

and then require that both $\|\nabla \tau\|_{L^p_{\delta-2}}^2$ and $\|((n-1)\tau+|\theta_-|/2)\|_{1-\frac{1}{p},p;\Sigma}$ be sufficiently small to obtain the inequalities in (5.12).

The final two theorems of this section provide us with a method to construct global subsolutions $\phi_- \leq \phi_+$, where ϕ_+ is any of the supersolutions constructed in Theorems 5.2, 5.3, or 5.5.

Theorem 5.6. (Global Subsolution for $g \in \mathcal{Y}^+$) Suppose that (\mathcal{M}, g) is asymptotically Euclidean of class $W^{2,p}_{\gamma}$, with n < p and $\gamma \in (2-n,0)$. Additionally assume that $\mathcal{Y}_g > 0$, $2-n < \delta < \gamma/2$, $\tau \in W^{1,p}_{\delta-1}$, $\rho \in L^p_{\gamma-2}$, $\sigma \in W^{1,2p}_{\gamma-1}$, $J \in \mathbf{L}^p_{\delta-2}$, $\theta_- \in W^{1-\frac{1}{p},p}(\Sigma)$, and $((n-1)\tau + |\theta_-|) > 0$ on Σ . Then there exists a subsolution $\phi_- > 0$ to the Hamiltonian constraint with boundary conditions (1.6)-(1.8) such that $\phi_- - \alpha\omega \in W^{2,p}_{\gamma}$ for $\alpha > 0$ sufficiently small.

Proof. Because $g \in \mathcal{Y}^+$, there exists $u \in W^{2,p}$ which solves

$$-\Delta u + (c_n R + b_n \tau^2) u = -\omega (c_n R + b_n \tau^2),$$

$$\partial_{\nu} u + \left(d_n H + d_n \tau + \frac{d_n}{(n-1)} |\theta_-| \right) u = -\omega \left(d_n H + d_n \tau + \frac{d_n}{(n-1)} |\theta_-| \right).$$
(5.13)

Set $\phi_- = \alpha(u + \omega)$, where $\alpha > 0$ will be determined. We observe that

$$-\Delta \phi_{-} + (c_{n}R + b_{n}\tau^{2})\phi_{-} = 0,$$

$$\partial_{\nu}\phi_{-} + \left(d_{n}H + d_{n}\tau + \frac{d_{n}}{(n-1)}|\theta_{-}|\right)\phi_{-} = 0,$$

and by the maximum the principles A.1 and A.2, $\phi_- > l > 0$ given that $\phi_- \to A_i > 0$ on each end.

We claim that for α sufficiently small, $\phi_- = \alpha \psi$ is a global subsolution. Suppose that $\phi \geqslant \phi_{-}$. Then we have

$$\mathcal{N}(\phi_-, S(\phi))$$

$$\leqslant \left(\begin{array}{c} b_n(\alpha^{N-1}(u+\omega)^{N-1} - \alpha(u+\omega))\tau^2 - c_n|\sigma + \mathcal{L}W|^2\phi_-^{-N-1} - c_n\rho\phi_-^{-\frac{N}{2}} \\ \left(\alpha^{\frac{N}{2}}(u+\omega)^{\frac{N}{2}} - \alpha(u+\omega)\right)\left(d_n\tau + \frac{d_n}{(n-1)}|\theta_-|\right) \end{array}\right),$$

where we have used the fact that $S(\nu,\nu)=((n-1)\tau+|\theta_-|/2)\psi^N>0$. We observe that if we take α sufficiently small, both expressions in the above array will be nonpositive.

Theorem 5.7. (Global Subsolution for bounded R and H) Let the assumptions of Lemma 5.4 hold along with additional assumption that $S(\nu, \nu) \ge 0$. Then there exists a solution u to

$$-\Delta u + c_n R u + b_n \tau^2 u^5 = 0 \quad \text{on } \mathcal{M},$$

$$\partial_{\nu} u + d_n H u + \left(d_n \tau + \frac{d_n}{n-1} |\theta_-| \right) u^{\frac{N}{2}} = 0 \quad \text{on } \Sigma,$$
(5.14)

such that $u-\omega\in W^{2,p}_{\gamma}$. Moreover, for any $\alpha\in(0,1)$ the function $\phi_-=\alpha u$ will be a global subsolution to the Hamiltonian constraint with boundary conditions (1.6)-(1.8).

Proof. We observe that $u_{-} \equiv 0$ is a subsolution of (5.14) and $u_{+} \equiv \beta \geq 1$ is a supersolution given the assumptions on R and H. Let ω have asymptotic limits $A_i \in (0, \infty)$. By Theorem A.4 we can choose β large enough so that Eq. (5.14) has a solution u such that $u-\omega\in W^{2,p}$. By construction, $u\geqslant 0$. We note that if $u(x_0)=0$ for some $x_0\in\mathcal{M}$, then x_0 will be a minimum of u. Both u and 0 satisfy the elliptic equation

$$\begin{split} -\Delta v + (c_n R + b_n \tau u^{N-2}) v &= 0 \quad \text{on } \mathcal{M}, \\ \partial_\nu v + (d_n H + \left(d_n \tau + \frac{d_n}{n-1} |\theta_-|\right) u^{\frac{N-2}{2}}) v &= 0 \quad \text{on } \Sigma, \end{split}$$

and u and the zero function will coincide up to first order at x_0 . Alexandrov's Theorem (cf. [2]) implies $u \equiv 0$, which contradicts the fact that $u \to A_i > 0$ on each E_i . So u > 0

Setting $\phi_- = \alpha u$ for $\alpha \in (0,1)$, we calculate $\mathcal{N}(\phi_-, S(\phi))$ for $\phi \geqslant \phi_-$:

$$\mathcal{N}(\phi_{-}, S(\phi)) = \begin{pmatrix} -c_{n} | \sigma + \mathcal{L}W |^{2} \phi_{-}^{-N-1} - c_{n} \rho \phi_{-}^{-\frac{N}{2}} \\ -\frac{d_{n}}{n-1} S(\nu, \nu) \phi_{-}^{-\frac{N}{2}} \end{pmatrix} \leqslant 0.$$

Therefore $\phi_- = \alpha u$ is a global subsolution to the Hamiltonian constraint.

Given $\omega_1 \in \mathcal{H}$ which tends to positive values on each end E_i , for arbitrarily small $\alpha>0$ we may obtain a positive subsolution ϕ_- such that $\phi_--\alpha\omega_1\in W^{2,p}_\gamma$. Similarly, given $\omega_2\in\mathcal{H}$ which tends to positive values on each end, for $\beta>0$ sufficiently small there exists a positive supersolution ϕ_+ such that $\phi_+-\beta\omega_2\in W^{2,p}_\delta$. By choosing $\alpha\ll\beta$, we can ensure that $\alpha\omega_1$ is asymptotically bounded by $\beta\omega_2$ and that $\phi_-\leqslant\phi_+$. Now that we have constructed barriers for the Hamiltonian constraint with the specified boundary conditions, we are ready to prove Theorem 3.1.

6. Non-CMC Solutions: Fixed Point Argument

Given a set of global barriers $\phi_- \leqslant \phi_+$, which we derived in Section 5, Theorem 3.1 will follow by using a variation of the fixed point argument first developed in [7]. The following argument closely follows the work done in [5], where the authors extended the argument in [7] to AE manifolds with no boundary. We slightly modify this fixed point argument to include our boundary problem.

Before we prove our fixed point theorem, we first discuss the properites of the solution map of the Hamiltonian constraint with the associated marginally trapped surface boundary conditions. In particular, we show that this map is well-defined up to the asymptotic limit of the solution, and then show that it is continuous.

Let $W \in W^{2,p}_{\delta}$ be a given vector field with $2-n < \delta \leq \gamma/2 < 0$, and let $\phi_- \leqslant \phi_+$ be sub-and supersolutions of $\mathcal{N}(\phi,W)$. For a given k-tuple A_1, \cdots, A_k of positive, real numbers, let $\omega \in \mathcal{H}$ be the associated harmonic function. By Theorems A.4 and A.5, for a given W, sub- and supersolutions ϕ_- and ϕ_+ , and ω that is asymptotically bounded by ϕ_- and ϕ_+ , there exists a unique solution to $\mathcal{N}(\phi,W)=0$ such that $\phi-\omega \in W^{2,p}_{\gamma}$. Therefore, for given a $W,\phi_- \leq \phi_+$, and ω , we define $T(W)=\phi$ to be the solution map giving this unique solution.

Given that T is used to construct our fixed point map for the Schauder Theorem, we require that the T(W) be a continuous mapping. We note that $\mathcal{G}(\phi)=i(T(\mathcal{S}(\phi)))$, where $i:\mathcal{H}+W^{2,p}_{\gamma}\to C^0$ is the compact embedding (2.4) and \mathcal{S} is the continuous solution map of the momentum constraint. Therefore the continuity of T will imply the continuity of T. We set T0 we set T1 and define T2 and define T3 will imply the continuity of T4. Then for fixed data T4 data T5 for a given 2-tensor T6. That is, T7 gives the solution of

$$-\Delta \phi + c_n R \phi + b_n \tau^2 \phi^{N-1} - c_n |\beta|^2 \phi^{-N-1} - c_n \rho \phi^{-\frac{N}{2}} = 0 \text{ on } \mathcal{M},$$

$$\partial_{\nu} \phi + d_n H \phi + \left(d_n \tau - \frac{d_n}{n-1} \theta_- \right) \phi^{\frac{N}{2}} - \frac{d_n}{n-1} \beta(\nu, \nu) \phi^{-\frac{N}{2}} = 0 \text{ on } \Sigma.$$

To prove the continuity of T it is sufficient to prove the continuity of \mathcal{L} in β . The proof is based on the Implicit Function Theorem argument developed in [10].

Proposition 6.1. Suppose (\mathcal{M},g) is asymptotically Euclidean of class $W^{2,p}_{\gamma}$, with $\gamma \in (2-n,0)$ and $2 > \frac{n}{p}$. Additionally assume that $\tau \in W^{1,p}_{\gamma/2-1}$, $\rho \in L^p_{\gamma-2}$, $\theta_- \in W^{1-\frac{1}{p},p}(\Sigma)$, and $\beta \in W^{1,2p}_{\gamma/2-1}$. If $((n-1)\tau + |\theta_-|) \geqslant 0$ and $\beta_0(\nu,\nu) \geqslant 0$, then \mathcal{L} is a C^1 map from $W^{1,2p}_{\gamma/2-1}$ to $W^{2,p}_{\gamma}$.

Proof. As in [10], we exploit the conformal covariance of the Lichnerowicz equation. Let $\hat{g} = \phi^{N-2}g$ and $\hat{\mathcal{L}}$ be the solution map associated with \hat{g} . By the conformal covariance of the boundary problem demonstrated in [8], we have that

$$\hat{\mathcal{L}}(\hat{\beta}) = \phi^{-1} \mathcal{L}(\beta) = 1$$
, where $\hat{\beta} = \phi^{-2N} \beta$.

Therefore it suffices to demonstrate the continuity of \mathcal{L} in a neighborhood of β_0 such that $\mathcal{L}(\beta_0) = 1$, and we may drop the hat notation.

Define

$$\mathcal{F}(\phi,\beta) = \begin{bmatrix} -\Delta\phi + c_n R\phi + b_n \tau^2 \phi^{N-1} - c_n |\beta|^2 \phi^{-N-1} - c_n \rho \phi^{-\frac{N}{2}} \\ \partial_{\nu}\phi + d_n H\phi + \left(d_n \tau + \frac{d_n}{(n-1)} |\theta_-| \right) \phi^{\frac{N}{2}} - \frac{d_n}{(n-1)} \beta(\nu,\nu) \phi^{-\frac{N}{2}} \end{bmatrix}.$$
(6.1)

It is clear that $\mathcal{F}(\mathcal{L}(\beta), \beta) = 0$, and a standard computation shows that the Gateaux derivative is given by

$$\mathcal{F}'_{\phi,\beta}(h,K) = \begin{pmatrix} -\Delta h + \alpha_1(\phi,\beta)h - 2c_n\phi^{-N-1}\beta \cdot K \\ \partial_{\nu}h + \alpha_2(\phi,\beta)h - \frac{d_n}{(n-1)}\phi^{-\frac{N}{2}}K(\nu,\nu) \end{pmatrix}, \tag{6.2}$$

where

$$\alpha_1(\phi,\beta) = c_n R + (N-1)b_n \tau^2 \phi^{N-2} + (N+1)c_n |\beta|^2 \phi^{-N-2} + \frac{c_n N}{2} \rho \phi^{-\frac{N}{2}-1},$$

and

$$\alpha_2(\phi,\beta) = d_n H + \frac{N}{2} \phi^{\frac{N}{2}-1} \left(d_n \tau + \frac{d_n}{(n-1)} |\theta_-| \right) + \frac{N d_n}{2(n-1)} \phi^{-\frac{N}{2}-1} \beta(\nu,\nu).$$

The multiplication properties of weighted Sobolev spaces imply that the operator \mathcal{F}' is continuous in ϕ and β . We have

$$\mathcal{F}'_{1,\beta_0}(h,0) = \begin{pmatrix} -\Delta h + \left(c_n R + (N-1)b_n \tau^2 + (N+1)c_n |\beta_0|^2 + \frac{c_n N}{2}\rho\right)h \\ \partial_{\nu} h + \left(d_n H + \frac{N}{2}\left(d_n \tau + \frac{d_n}{(n-1)}|\theta_-|\right) + \frac{Nd_n}{2(n-1)}\beta_0(\nu,\nu)\right)h \end{pmatrix},$$

and given that $\mathcal{F}(1,\beta_0) = \mathbf{0}$,

$$c_n R + b_n \tau^2 - c_n |\beta_0|^2 - c_n \rho = 0,$$

 $d_n H + \left(d_n \tau + \frac{d_n}{(n-1)} |\theta_-| \right) - \frac{d_n}{(n-1)} \beta_0(\nu, \nu) = 0.$

This implies that

$$\mathcal{F}'_{1,\beta_0}(h,0) = \begin{pmatrix} -\Delta h + \left((N-2)b_n \tau^2 + (N+2)c_n |\beta_0|^2 + \frac{N+2}{2}c_n \rho \right) h \\ \partial_{\nu} h + \left(\frac{N-2}{2} \left(d_n \tau + \frac{d_n}{(n-1)} |\theta_-| \right) + \frac{(N+2)}{2} \frac{d_n}{(n-1)} \beta_0(\nu, \nu) \right) h \end{pmatrix}. \quad (6.3)$$

The assumptions $\beta_0(\nu,\nu)\geqslant 0$ and $((n-1)\tau+|\theta_-|)\geqslant 0$ imply that $\mathcal{F}'_{1,\beta_0}:W^{2,p}_{\gamma}\to L^p_{\gamma-2}$ is an isomorphism, and the Implicit Function Theorem implies that \mathcal{L} is continuous in a neighborhood of β_0 .

Now that we have established existence of global barriers and showed that \mathcal{G} is continuous, we are ready to prove Theorems 3.1 and 3.2.

Proof of Theorem 3.1 Let \mathcal{C}^0_+ denote the set of strictly positive bounded functions on \mathcal{M} . If $\phi \in \mathcal{C}^0_+$, then by Proposition 4.1, the vector field $W = \mathcal{S}(\phi) \in W^{2,p}_\delta$ given by the solution map of the momentum constraint with boundary conditions (1.6) is well-defined. By the remarks preceding Proposition 6.1 and Theorems A.4 and A.5, given $W \in W^{2,p}_\delta$, sub-and super-solutions $\phi_- \leqslant \phi_+$, and a harmonic function ω as in Proposition A.3 that is asymptotically bounded by ϕ_- and ϕ_+ , the solution map $T(W) = \varphi$ is well-defined and continuous.

Let $\omega_1, \omega_2 \in \mathcal{H}$ tend to positive real numbers on each end and suppose that ω_2 is asymptotically bounded above by ω_1 . Let ϕ_+ be the global supersolution obtained from

either Theorem 5.2 or 5.3, where $\phi_+ - \beta \omega_1 \in W_{\gamma}^{2,p}$. Note that we use the supersolution from Theorem 5.2 if we wish to solve the coupled constraints with boundary conditions (1.6)-(1.8) with arbitrary ψ . If we wish to obtain a solution ϕ satisfying the marginally trapped surface condition $\phi \leqslant \psi$, then we use the supersolution from Theorem 5.3. Let $\phi_- \leqslant \phi_+$ be the global subsolution obtained from Theorem 5.6, where $\phi_- - \alpha \omega_2 \in W_{\gamma}^{2,p}$. Let $\omega \in \mathcal{H}$ be asymptotically bounded by $\phi_- \leqslant \phi_+$. For this choice of sub-and supersolutions and ω we may apply Theorem A.4 to obtain $\varphi = T(W)$. Following the proof of Theorem A.4, $\varphi = \omega + \hat{\varphi} \in \mathcal{H} + W_{\gamma}^{2,p}$. Let i denote the compact inclusion $\mathcal{H} + W_{\delta}^{2,p} \hookrightarrow \mathcal{C}^0$. A solution (ϕ, W) to (5.1) then corresponds to a fixed point of the mapping $\mathcal{G}(\phi) = i(T(\mathcal{S}(\phi)))$, which is a continuous, compact mapping.

Define the bounded convex set $\mathcal{A} := \{ \phi \in \mathcal{C}^0_+ : \phi_- \leqslant \phi \leqslant \phi_+ \}$. By construction, \mathcal{G} maps \mathcal{A} to itself. Moreover, \mathcal{A} is closed, bounded, and convex. Therefore the Schauder fixed point theorem implies that \mathcal{A} contains a fixed point ϕ of \mathcal{G} . Standard estimates imply that ϕ and $W(\phi)$ both have the desired regularity. \square

Proof of Theorem 3.3 The proof if the same as the proof of Theorem 3.1 except for the barriers used. Given $A_i \in [1, \infty)$, let $\omega \in \mathcal{H}$ be the associated harmonic function. Choose $\omega_1, \omega_2 \in \mathcal{H}$ such that $\omega_1 \leqslant \omega_2$ and ω is asymptotically bounded by ω_1 and ω_2 . By Theorem 5.5 there exists a global supersolution $\phi_+ = \phi_A$ to the Hamiltonian constraint with boundary conditions (1.6)-(1.8) such that $\phi_+ - \omega_2 \in W_\gamma^{2,p}$. By setting $\psi = \phi_+$ we can impose the marginally trapped surface condition $\phi \leqslant \phi_+$. By Theorem 5.7 we obtain a global subsolution $\phi_- \leqslant \phi_+$ such that $\phi_- - \alpha\omega_1 \in W_\gamma^{2,p}$ for $\alpha \in (0,1)$. By construction, the function ω is asymptotically bounded by $\phi_- \leqslant \phi_+$, and as in the proof of Theorem 3.1 we apply Theorem A.4 to obtain a solution φ such that $\varphi - \omega \in W_\gamma^{2,p}$. Therefore, for given asymptotic limits $A_i \in [1, \infty)$, the solution map $\varphi = T(W(\phi))$ is well-defined for $\phi_- \leqslant \phi \leqslant \phi_+$. The rest of the proof follows from the arguments made in the proof of Theorem 3.1.

7. NEAR-CMC SOLUTIONS: AN IMPLICIT FUNCTION THEOREM ARGUMENT

In this section, we provide an alternative approach to obtain solutions to the conformal equations satisfying the marginally trapped surface boundary conditions. This approach is based on the Implicit Function Theorem argument given in [2], and therefore requires that $\|\tau\|_{W^{1,p}_s}$ be sufficiently small.

We first recall the Implicit Function Theorem. Suppose that U and V are open subsets of Banach spaces X and Y and \mathcal{F} is a C^1 mapping from $U \times V$ into a Banach space Z:

$$\mathcal{F}: X \times Y \to Z$$
.

The Implicit function theorem states that if $\mathcal{F}_y(x_0,y_0)$ is invertible at some solution of $\mathcal{F}(x_0,y_0)=0$, then there exists a neighborhood $U'\times V'\subset U\times V$ of (x_0,y_0) such that for each $x\in U'$, there exists a unique $y\in V'$ such that $\mathcal{F}(x,y)=0$. That is, there exists an invertible function $\rho:V'\to U'$ such that all solutions to $\mathcal{F}(x,y)=0$ in $U'\times V'$ are of the form $F(\rho(y),y)=0$. Moreover, if \mathcal{F} is C^1 then ρ is C^1 .

Given a k-tuple A_1, \dots, A_k of positive numbers in \mathbb{R} , let $\omega \in \mathcal{H}$ be the associated harmonic function. Suppose that $2 - n < \delta \le \gamma/2 < 0$. As in [2], define the variables

$$x = (\tau, J) \in X = W_{\delta-1}^{1,p} \times L_{\delta-2}^{p}(T\mathcal{M}),$$

$$y = (\phi - \omega, W) \in Y = (W_{\gamma}^{2,p} \times W_{\delta}^{2,p}(T\mathcal{M})) \cap \{\phi > 0\},$$

$$Z = L_{\gamma-2}^{p} \times W^{1-\frac{1}{p},p}(\Sigma) \times L_{\delta-2}^{p}(T\mathcal{M}) \times W^{1-\frac{1}{p},p}(T\Sigma).$$
(7.1)

Let

$$\mathcal{F}(x,y) = \begin{bmatrix} -\Delta\phi + c_n R\phi + b_n \tau^2 \phi^{N-1} - c_n |\sigma + \mathcal{L}W|^2 \phi^{-N-1} - c_n \rho \phi^{-\frac{N}{2}} \\ \partial_{\nu}\phi + d_n H\phi \\ \Delta_{\mathbb{L}}W + \frac{n-1}{n} \nabla \tau \phi^N + J \\ \mathcal{L}W(\nu, \cdot) - \mathbf{V}(\phi, \tau) \end{bmatrix}, (7.2)$$

where $\mathbf{V}(\phi,\tau)(\nu)=(n-1)\tau\phi^N-\sigma(\nu,\nu)$, which implies that $S(\nu,\nu)=(n-1)\tau\phi^N$. We observe that solutions to $\mathcal{F}(x,y)=0$ represent solutions to the coupled system (1.4)-(1.5) with boundary conditions (1.6)-(1.8) when $|\theta_-|=0$ and $\psi=\phi$. These solutions will satisfy the marginally trapped surface conditions if $\tau\geqslant 0$ on Σ .

Fix $\sigma \in W^{1,2p}_{\gamma/2-1}$ and $\rho \in L^p_{\gamma-2}$. In order to apply the Implicit Function Theorem to (7.2), we require that $\mathbf{V} = \mathbf{V}(\phi, \tau)$ be a C^1 vector field in (ϕ, τ) . We construct such a \mathbf{V} in the following proposition.

Proposition 7.1. Suppose that (\mathcal{M},g) is an n-dimensional asymptotically Euclidean manifold of class $W^{2,p}_{\gamma}$ with $\gamma \in (2-n,0)$ and compact boundary Σ . If $2-n < \delta < \gamma/2$, $\tau \in W^{1,p}_{\delta-1}$, and $\phi \in C^0$, then there exists a vector field $\mathbf{V}(\phi,\tau) \in \mathbf{W}^{1,p}$ that is C^1 in ϕ and τ . Moreover, \mathbf{V} satisfies

$$\mathbf{V}(\nu) = (n-1)\tau\phi^N - \sigma(\nu,\nu). \tag{7.3}$$

Proof. For every point $p \in \Sigma$ there is a neighborhood U and a coordinate map Ψ such that $\Psi(p) \in \{\mathbf{x} \in \mathbb{R}^n \mid x_n = 0\}$ and $V = \Psi(U) \subseteq \{\mathbf{x} \in \mathbb{R}^n \mid x_n \geqslant 0\}$. There exists a radius R > 0 such that $B_R(\Psi(p)) \cap \{\mathbf{x} : x_n \geqslant 0\} \subset V$. Let

$$A = \Psi^{-1}(B_R(\Psi(p)) \cap {\mathbf{x} : x_n \ge 0}).$$

On V, we define the constant vector field $\mathbf{W}(\mathbf{x}) = (0, \cdots, 0, -1)$, and then consider the pullback $\mathbf{X} = \Psi^*(\mathbf{W})$ on A. By construction, $\mathbf{X} = \nu$ on $\Sigma \cap A$. The compactness of Σ implies that there exists some collection of $p_i \in \Sigma$ such that the associated sets A_i as above determine a finite covering of Σ for $1 \leq i \leq M$. Let \mathbf{X}_i be the associated local vector fields defined on A_i , and let A_0 be an open set such that $\bigcup_{i=0}^M A_i = \mathcal{M}$. Finally, let χ_i be a partition of unity subordinate to the covering $\{A_i\}$. Setting

$$\mathbf{V} = \sum_{i=0}^{M} \chi_i((n-1)\tau\phi^N - \sigma(\nu,\nu))\mathbf{X}_i = ((n-1)\tau\phi^N - \sigma(\nu,\nu))\mathbf{X},$$

where $\mathbf{X}_0 = \mathbf{0}$ on A_0 , it is clear that $\mathbf{V} = ((n-1)\tau\phi^N - \sigma(\nu,\nu))\nu$ on Σ given that $\mathbf{X} = \nu$ on Σ . By construction, $\mathbf{V} \in \mathbf{W}^{1,p}$ given the regularity assumptions on τ , σ and ϕ and the fact that \mathbf{V} vanishes outside of a neighborhood of Σ . Clearly \mathbf{V} will be C^1 in ϕ and τ .

The properties of the vector ${\bf V}$ constructed in Proposition 7.1 and the multiplication properties of weighted Sobolev spaces imply that ${\cal F}(x,y): X\times Y\to Z$ will be C^1 as long as p>n and $\delta\in(2-n,0)$ (cf. [1, 2]). Letting ${\bf V}=((n-1)\tau\phi^N-\sigma(\nu,\nu)){\bf X}$ as in the proof of Proposition 7.1, the partial derivative for a given (x,y)

$$\mathcal{F}'_y(x,y): Y \to Z,$$

 $(h,\beta): \to \mathcal{F}'_y(x,y)(h,\beta),$

is given by

$$\mathcal{F}'_{y}(x,y)(h,\beta) = \begin{pmatrix} -\Delta h + \alpha(\phi,W)h - 2c_{n}\phi^{-N-1}(\sigma + \mathcal{L}W) \cdot \mathcal{L}\beta \\ \partial_{\nu}h + d_{n}Hh \\ \Delta_{\mathbb{L}}\beta + N\frac{(n-1)}{n}\nabla\tau\phi^{N-1}h \\ \mathcal{L}\beta(\nu,\cdot) - N(n-1)\tau\phi^{N-1}\mathbf{X}h \end{pmatrix}, \tag{7.4}$$

where

$$\alpha(\phi, W) = c_n R + (N-1)b_n \tau^2 \phi^{N-2} + (N+1)c_n |\sigma + \mathcal{L}W|^2 \phi^{-N-2} + \frac{c_n N}{2} \rho \phi^{-\frac{N}{2}-1}$$

and $X = \nu$ is a smooth vector field on \mathcal{M} vanishing in a neighborhood of Σ .

Theorem 7.2. Suppose that (\mathcal{M},g) is an n-dimensional asymptotically Euclidean manifold of class $W^{2,p}_{\gamma}$, where $\gamma \in (2-n,0)$ and p > n. Assume that $\sigma \in W^{1,2p}_{\gamma/2-1}$ and $\rho \in L^p_{\gamma-2}$ are given. Suppose that $\mathcal{F}(x,y) = 0$ has a solution when $y_0 = (\tau_0, \mathbf{J}_0) = (0,\mathbf{0})$, and denote this solution by $x_0 = (\phi_0,W_0)$. If $\alpha(\phi_0,W_0) \geqslant 0$ and $H \geqslant 0$, then there exists a neighborhood U of (τ_0,\mathbf{J}_0) in X such that the coupled constraints with boundary conditions (1.6)-(1.7) have a unique solution (ϕ,W) , $\phi > 0$, $(\phi - \omega,W) \in Y$.

Proof. We calculate

$$\mathcal{F}'_{y}(x_{0}, y_{0})(h, \beta) = \begin{pmatrix} -\Delta h + \alpha(\phi_{0}, W_{0})h - 2c_{n}\phi_{0}^{-N-1}(\sigma + \mathcal{L}W_{0}) \cdot \mathcal{L}\beta \\ \partial_{\nu}h + d_{n}Hh \\ \Delta_{\mathbb{L}}\beta \\ \mathcal{L}\beta(\nu, \cdot) \end{pmatrix}. \tag{7.5}$$

By Proposition 4.1 the operator $\mathcal{F}'_y(x_0,y_0):Y\to Z$ is invertible. Therefore the Implicit Function Theorem implies the result.

Corollary 7.3. Suppose that (\mathcal{M},g) is an n-dimensional asymptotically Euclidean manifold of class $W^{2,p}_{\gamma}$, where $\gamma \in (2-n,0)$ and p > n. Assume that $\sigma \in W^{1,2p}_{\gamma/2-1}$ and $\rho \in L^p_{\gamma-2}$ are given. Let $x_0 = (\phi_0,W_0)$ denote the solution to $\mathcal{F}(x,y) = 0$ when $y_0 = (\tau_0,\mathbf{J}_0) = (0,0)$. Then there exists a neighborhood of (ϕ_0,W_0) in which solutions to $\mathcal{F}(x,y) = 0$ exist and are unique. In particular, there exist unique solutions (ϕ,W) in this neighborhood where $(\phi-\omega,W) \in Y$ satisfies the marginally trapped surface conditions.

Proof. The existence and uniqueness of (ϕ_0, W_0) such that $\phi_0 - \omega \in W^{2,p}_{\gamma}$ follows from Section 8 in [2]. By Proposition 4.1 the assumption that $g \in \mathcal{Y}^+$ also implies that $\mathcal{F}'_y(x_0, y_0)$ in (7.5) is invertible. Therefore, we may apply the Implicit Function Theorem to uniquely parametrize $(\phi, W) \in X$ in terms of (τ, J) in a neighborhood of (ϕ_0, W_0) . Those solutions in a neighborhood of (ϕ_0, W_0) which correspond to $\tau \geqslant 0$ will satisfy the marginally trapped surface conditions.

APPENDIX A. SOLUTIONS TO SEMILINEAR, BOUNDARY VALUE PROBLEMS

Suppose that (\mathcal{M},g) is an n-dimensional, asymptotically Euclidean manifold with boundary Σ of class $W^{2,p}_{\gamma}$, with $\gamma \in (2-n,0)$ and p>n. Denote the ends of \mathcal{M} by E_i for $1 \leq i \leq m$. Here we investigate the existence of solutions to the semilinear, Robin problem

$$-\Delta u = f_1(x, u) \quad \text{on } \mathcal{M},$$

$$\partial_{\nu} u = f_2(x, u) \quad \text{on } \partial \mathcal{M}.$$
(A.1)

The functions $f_i(x,y): \mathcal{M} \times I_i \to \mathbb{R}$ for $i \in \{1,2\}$ are of the form

$$f_i(x,y) = \sum_{j=1}^{N_i} a_{ij}(x)b_{ij}(y),$$

where each $b_{ij}(y)$ is a smooth function on $I_i \subset \mathbb{R}$, and $a_{ij}(x) \in L^p_{\gamma-2}$.

In order to develop an iterative method which solves (A.1), we will require the following version of the weak maximum principle.

Lemma A.1. Suppose that (\mathcal{M}, g) satisfies the assumptions above, and that $V(x) \in L^p_{\gamma-2}$ and $\mu(x) \in W^{1-\frac{1}{p},p}(x)$ are nonnegative. If

$$-\Delta u + V(x)u \ge 0 \quad on \,\mathcal{M},$$
$$\partial_{\nu} u + \mu(x)u \ge 0 \quad on \,\partial\mathcal{M},$$

and $u \to A_i \ge 0$ on each end E_i , then $u \ge 0$.

Proof. Let w = -u. Given that $u \to A_i \ge 0$ on each end, the function $v = (w - \epsilon)^+$ has compact support. By Sobolev embedding $v \in W^{1,2}$, and $wv \ge 0$. We have

$$\|\nabla v\|_{L^{2}(\mathcal{M})}^{2} = \int_{\mathcal{M}} \nabla w \cdot \nabla v \, dV = -\int_{\mathcal{M}} (\Delta w) v \, dV + \int_{\partial \mathcal{M}} (\partial_{\nu} w) v \, dA$$

$$\leq -\int_{\mathcal{M}} V(x) w v \, dV - \int_{\partial \mathcal{M}} \mu(x) w v \, dA \leq 0.$$

Therefore $v \equiv 0$ and $u \ge -\epsilon$ on \mathcal{M} . Letting $\epsilon \to 0$ we have that $u \ge 0$.

We also require a version of the strong maximum taken from [9]. For completeness, we state it here without proof.

Lemma A.2. Suppose that (\mathcal{M},g) satisfies the assumptions above and $V(x) \in L^p_{\gamma-2}$ and $\mu(x) \in W^{1-\frac{1}{p},p}(x)$. Suppose $u(x) \in W^{2,p}_{\gamma}$ is nonnegative and

$$-\Delta u + V(x)u \ge 0 \quad \text{on } \mathcal{M},$$
$$\partial_{\nu} u + \mu(x)u \ge 0 \quad \text{on } \partial \mathcal{M}.$$

If u(x) = 0 for some $x \in \mathcal{M}$, then u vanishes identically.

In the following Lemma we construct an auxiliary, harmonic function which allows us to freely specify the asymptotic limit A_j on each end E_j of the solution to (A.1). This argument is a modification of an argument given in [5].

Proposition A.3. Suppose \mathcal{M} has ends E_1, \dots, E_k , and let $A_j \in (-\infty, \infty)$ for $1 \leq j \leq k$. Then there exists a unique function ω solving

$$-\Delta\omega = 0 \quad on \ \mathcal{M},$$
$$\partial_{\nu}\omega = 0 \quad on \ \Sigma,$$

which tends to A_j on each end E_j . Moreover, $\min A_j \leq \omega \leq \max A_j$.

Proof. Let $\omega_1 = \sum \chi_j A_j$, where each χ_j is a cutoff function which equals 1 on the end E_j and is zero outside a neighborhood of this end. Then $\Delta\omega_1 \in L^p_{\delta-2}$ and $\partial_\nu\omega_1 \in W^{1-\frac{1}{p},p}$, so there exists a function $\omega_2 \in W^{2,p}_{\delta}$ such that $\Delta\omega_2 = \Delta\omega_1$ and $\partial_\nu\omega_2 = \partial_\nu\omega_1$. Therefore $\omega = \omega_1 - \omega_2$ has the desired properties. The fact that ω is unique follows by assuming that two such functions exist and showing that the difference must be identically zero.

To show that $\min A_j \leq \omega \leq \max A_j$, we first pick $\epsilon < \min A_j$ and define $v = (\omega - \epsilon)^- \in W^{1,2}$. Given the asymptotic behavior of ω , the function v has compact support. Therefore,

$$\|\nabla v\|_{L^2}^2 = \int_{\mathcal{M}} \nabla \omega \nabla v = 0,$$

which implies that $v \equiv 0$. Therefore $\omega \geq \epsilon$, and by letting $\epsilon \to \min A_j$ we have $\omega \geq \min A_j$. To show that $\omega \leq \max A_j$, we make a similar argument using $v = (\omega - \epsilon)^+$ for $\epsilon > \max A_j$ and let $\epsilon \to \max A_j$.

With these results in hand, we are now ready to address the existence of solutions to (A.1). The proof of the following theorem provides an iterative method to construct solutions to this problem given sub-and supersolutions, where we recall that a sub-solution ϕ_{-} for (A.1) satisfies

$$-\Delta \phi_{-} \le f_1(x, \phi_{-}),$$

$$\partial_{\nu} \phi_{-} \le f_2(x, \phi_{-}),$$

and a supersolutions ϕ_+ satisfies

$$-\Delta \phi_{+} \ge f(x, \phi_{+}),$$

$$\partial_{\nu} \phi_{+} \ge f_{2}(x, \phi_{+}).$$

Our argument is based on the construction in [2].

Theorem A.4. Suppose that (A.1) admits a subsolution and supersolution $\phi_-, \phi_+ \in W^{2,p}_{\gamma}$, and assume that

$$\ell < \phi_{-} < \phi_{+} < m, \quad [l, m] \subset I_1 \cap I_2,$$

and

$$\lim_{|x|\to\infty}\phi_-=\alpha,\quad \lim_{|x|\to\infty}\phi_+=\beta.$$

Let ω be as in Proposition A.3, where each A_j satisfies $\alpha \leq A_j \leq \beta$. Then Eq.(A.1) admits a solution ϕ such that

$$\phi_- \le \phi \le \phi_+, \quad \phi - \omega \in W^{2,p}_{\gamma}.$$

Proof. As in [2], the proof is by induction starting with ϕ_- . Let $k_1 \in L^p_{\gamma-2}$ and $k_2 \in W^{1-\frac{1}{p},p}$ be positive functions such that

$$k_1(x) \ge \sup_{1 \le y \le m} \frac{\partial}{\partial y} f_1(x, y), \quad \text{and} \quad k_2(x) \ge \sup_{1 \le y \le m} \frac{\partial}{\partial y} f_2(x, y).$$

We recall that by Proposition A.3, $\min A_j \leq \omega \leq \max A_j$. Setting $\phi_1 = \omega + u_1$, where u_1 satisfies

$$-\Delta u_1 + k_1 u_1 = f_1(x, \phi_-) + k_1(\phi_- - \omega)$$

$$\partial_{\nu} u_1 + k_2 u_1 = f_2(x, \phi_-) + k_2(\phi_- - \omega),$$
(A.2)

we conclude that

$$-\Delta(\phi_1 - \phi_-) + k_1(\phi_1 - \phi_-) \ge 0$$
$$\partial_{\nu}(\phi_1 - \phi_-) + k_2(\phi_1 - \phi_-) \ge 0.$$

By assumption, $\phi_1 - \phi_-$ tends to $A_i - \alpha \ge 0$ on each end E_i and Lemma A.1 implies that

$$\phi_1 \geq \phi_-$$
 on \mathcal{M} .

Similarly,

$$-\Delta(\phi_{+} - \phi_{1}) + k_{1}(\phi_{+} - \phi_{1}) \ge f_{1}(x, \phi_{+}) - f_{1}(x, \phi_{-}) + k_{1}(\phi_{+} - \phi_{-}),$$

$$\partial_{\nu}(\phi_{+} - \phi_{1}) + k_{2}(\phi_{+} - \phi_{1}) \ge f_{2}(x, \phi_{+}) - f_{2}(x, \phi_{-}) + k_{2}(\phi_{+} - \phi_{-}).$$
(A.3)

By our choice of $k_i(x)$, the function

$$g_i(x,y) = f_i(x,y) + k_i(x)y \tag{A.4}$$

is monotonic increasing in the variable y. Therefore both equations in (A.3) are nonnegative. Because $\phi_+ - \phi_1$ tends to $\beta - A_i \ge 0$ on each end E_i , we may apply the maximum principle A.1 again to conclude that $\phi_1 \le \phi_+$.

We now define $\phi_n = \omega + u_n$ inductively by letting

$$-\Delta u_n + k_1 u_n = f_1(x, \phi_{n-1}) + k_1 u_{n-1}$$
$$\partial_{\nu} u_n + k_2 u_n = f_2(x, \phi_{n-1}) + k_2 u_{n-1}.$$

Standard elliptic theory implies that $u_n \in W^{2,p}_{\gamma}$ for each n and $u_n \to \omega$ on each end. Assume that for all $1 \le i \le k \le n-1$, ϕ_i and ϕ_k are defined as above and satisfy $\phi_- \le \phi_i \le \phi_k \le \phi_+$. Then we have

$$-\Delta(\phi_n - \phi_{n-1}) + k_1(\phi_n - \phi_{n-1}) = f_1(x, \phi_{n-1}) - f_1(x, \phi_{n-2}) + k_1(\phi_{n-1} - \phi_{n-2}) \ge 0,$$

$$\partial_{\nu}(\phi_n - \phi_{n-1}) + k_2(\phi_n - \phi_{n-1}) = f_2(x, \phi_{n-1}) - f_2(x, \phi_{n-2}) + k_2(\phi_{n-1} - \phi_{n-2}) \ge 0,$$

where the above inequalities follow from the inductive hypothesis and the fact that (A.4) is monotonic increasing. As $\phi_n - \phi_{n-1} \to 0$ on each end E_i , the maximum principle implies that $\phi_n \ge \phi_{n-1}$. Finally, an application of the maximum principle to

$$-\Delta(\phi_{+} - \phi_{n}) + k_{1}(\phi_{+} - \phi_{n}) \ge f_{1}(x, \phi_{+}) - f_{1}(x, \phi_{n-1}) + k_{1}(\phi_{+} - \phi_{n-1}) \ge 0,$$

$$\partial_{\nu}(\phi_{+} - \phi_{n}) + k_{2}(\phi_{+} - \phi_{n}) \ge f_{2}(x, \phi_{+}) - f_{2}(x, \phi_{n-1}) + k_{2}(\phi_{+} - \phi_{n-1}) \ge 0,$$

implies that $\phi_n \leq \phi_+$.

Therefore, the sequence of functions $\phi_n \in W^{2,p}_{\gamma}$ is monotonic increasing and bounded above by ϕ_+ . Thus the sequence converges to a function $\phi(x) = \omega + u(x)$, with $\phi_- \le \phi \le \phi_+$. A standard bootstrapping argument as in [2] then implies that $\phi(x)$ has the desired regularity.

Theorem A.5. Let (\mathcal{M},g) be an n-dimensional, asymptotically Euclidean manifold with boundary Σ of class $W^{2,p}_{\gamma}$, where $\gamma \in (2-n,0)$ and p > n. Let $\tau \in W^{1,p}_{\gamma-1}$, $\sigma \in L^{2p}_{\gamma-1}$, $\rho \in L^p_{\gamma-2}$, $\theta_- \in W^{1-\frac{1}{p},p}(\Sigma)$ be fixed data such that $((n-1)\tau+|\theta_-|)\geq 0$ on Σ , and suppose that $W\in W^{2,p}_{\delta}$ is given, where $2-n<\delta\leqslant \gamma/2$ and $S(\nu,\nu)\geqslant 0$. Additionally assume that $A_i\in (0,\infty)$ and that ω is the associated smooth, harmonic function such that $\omega\to A_i$ on each end E_i . Finally, assume that the Lichnerowicz equation (1.4) with boundary conditions (1.6) has a sub- and supersolution ϕ_- and ϕ_+ which asymptotically bound ω . Then there exists a unique solution $\phi>0$ to the Lichnerowicz equation (1.4) with boundary conditions (1.6) such that $\phi-\omega\in W^{2,p}_{\gamma}$.

Proof. The fact that a solution exists follows from Theorem A.4. To see that this solution is unique, suppose that ϕ_1 and ϕ_2 are both solutions. That is, suppose that for $i \in \{1, 2\}$

$$\begin{split} -\Delta\phi_i + c_nR\phi_i + b_n\tau^2\phi_i^{N-1} - c_n|\sigma + \mathcal{L}W|^2\phi_i^{-N-1} - c_n\rho\phi_i^{-\frac{N}{2}} &= 0 \quad \text{on } \mathcal{M}, \\ \partial_\nu\phi_i + d_nH\phi_i + \left(d_n\tau - \frac{d_n}{n-1}\theta_-\right)\phi_i^{\frac{N}{2}} - \frac{d_n}{n-1}S(\nu,\nu)\phi_i^{-\frac{N}{2}} &= 0 \quad \text{on } \Sigma. \end{split}$$

The conformal transformation properties of the scalar curvature and boundary, mean extrinsic curvature then imply that

$$c_n R(\phi_i^{N-2}g)\phi_i^{N-1} = -b_n \tau^2 \phi_i^{N-1} + c_n |\sigma + \mathcal{L}W|^2 \phi_i^{-N-1} + c_n \rho \phi_i^{-\frac{N}{2}},$$

$$d_n H(\phi_i^{N-2}g)\phi_i^{\frac{N}{2}} = -\left(d_n \tau - \frac{d_n}{n-1}\theta_-\right)\phi_i^{\frac{N}{2}} + \frac{d_n}{n-1}S(\nu,\nu)\phi_i^{-\frac{N}{2}},$$

where $R(\phi_i^{N-2}g)$ and $H(\phi_i^{N-2}g)$ denote the scalar curvature and boundary mean curvature with respect to the metric $\phi_i^{N-2}g$. Setting $u=\phi_1^{-1}\phi_2$, we clearly have that $u-1\in W^{2,p}_\gamma$. Moreover, the above equation implies that

$$-\Delta_{\phi_{1}^{N-2}g}u + c_{n}(-b_{n}\tau^{2} + c_{n}|\sigma + \mathcal{L}W|^{2}\phi_{1}^{-2N} + c_{n}\rho\phi_{1}^{-\frac{3N}{2}+1})u =$$

$$c_{n}(-b_{n}\tau^{2} + c_{n}|\sigma + \mathcal{L}W|^{2}\phi_{2}^{-2N} + c_{n}\rho\phi_{2}^{-\frac{3N}{2}+1})u^{N-1},$$

$$\partial_{\nu}u + d_{n}\left(-\left(d_{n}\tau - \frac{d_{n}}{n-1}\theta_{-}\right) + \frac{d_{n}}{n-1}S(\nu,\nu)\phi_{1}^{-N}\right)u =$$

$$d_{n}\left(-\left(d_{n}\tau - \frac{d_{n}}{n-1}\theta_{-}\right) + \frac{d_{n}}{n-1}S(\nu,\nu)\phi_{2}^{-N}\right)u^{\frac{N}{2}},$$

where ∂_{ν} is with respect to $\phi_1^{N-2}g$. We note that the above equations have the form

$$\begin{split} -\Delta_{\phi_1^{N-2}g}u + (a + b\phi_1^{-2N} + c\phi_1^{-\frac{3N}{2}+1})u &= (a + b\phi_2^{-2N} + c\phi_2^{-\frac{3N}{2}+1})u^{N-1} \quad \text{on } \mathcal{M}, \\ \partial_{\nu}u + (e + f\phi_1^{-N})u &= (e + f\phi_2^{-N})u^{\frac{N}{2}} \quad \text{on } \Sigma, \end{split}$$

where $a \le 0, b \ge 0, c \ge 0$ and $e \le 0, f \ge 0$. Rearranging these two equations, we obtain

$$-\Delta_{\phi_1^{N-2}g}u + \frac{1}{(u-1)} \left(a(1-u^{N-2}) + b\phi_1^{-2N}(1-u^{-N-2}) + c\phi_1^{-\frac{3N}{2}+1}(1-u^{-\frac{N}{2}-1}) \right) u(u-1) = 0,$$

$$\partial_{\nu}u + \frac{1}{(u-1)} \left(e(1-u^{\frac{N}{2}-1}) + f\phi_1^{-N}(1-u^{-\frac{N}{2}-1}) \right) u(u-1) = 0.$$

We observe that for m>0, $\frac{u^m-1}{(u-1)}>0$ given that u>0. This also implies that $\frac{(1-u^{-m})}{(u-1)}=u^{-m}\frac{u^m-1}{(u-1)}>0$. Given the assumptions on a,b,c,e and f, we conclude that

$$\begin{split} \frac{1}{(u-1)} \left(a(1-u^{N-2}) + b\phi_1^{-2N}(1-u^{-N-2}) + c\phi_1^{-\frac{3N}{2}+1}(1-u^{-\frac{N}{2}-1}) \right) u \geqslant 0 \text{ on } \mathcal{M}, \\ \frac{1}{(u-1)} \left(e(1-u^{\frac{N}{2}-1}) + f\phi_1^{-N}(1-u^{-\frac{N}{2}-1}) \right) u \geqslant 0 \text{ on } \Sigma. \end{split}$$

The fact that $u-1 \in W^{2,p}_{\gamma}$ and Lemma A.1 imply that $u-1 \geq 0$ on \mathcal{M} . Thus, $\phi_2 \geq \phi_1$. We may obtain the inequality $\phi_1 \geq \phi_2$ by reversing the roles of ϕ_1 and ϕ_2 in the above argument. Thus, $\phi_1 = \phi_2$ and the solution is unique.

REFERENCES

- [1] R. Bartnik. The mass of an asymptotically flat manifold. *Comm. Pure Appl. Math*, (39):661–693, 1986.
- [2] Y. Choquet-Bruhat, J. Isenberg, and J. York. Einstein constraint on asymptotically Euclidean manifolds. *Phys. Rev. D*, 61:084034, 2000.
- [3] S. Dain. Trapped surfaces as boundaries for the constraint equations. *Classical Quantum Gravity*, 21(2):555–573, 2004.

- [4] J. Dilts. A limit equation for the solvability of the constraint equations in the asymptotically flat case.
- [5] J. Dilts, J. Isenberg, R. Mazzeo, and C. Meier. Non-CMC solutions to the Einstein constraint equations on ayamptotically Euclidean manifolds. Class. Quantum Grav., (31):065001, 2014.
- [6] M. Holst, G. Nagy, and G. Tsogtgerel. Far-from-constant mean curvature solutions of Einstein's constraint equations with positive Yamabe metrics. Phys. Rev. Lett., 100(16):161101.1-161101.4, 2008. Available as arXiv:0802.1031 [gr-qc].
- [7] M. Holst, G. Nagy, and G. Tsogtgerel. Rough solutions of the Einstein constraints on closed manifolds without near-CMC conditions. Comm. Math. Phys., 288(2):547-613, 2009. Available as arXiv:0712.0798 [gr-qc].
- [8] M. Holst and G. Tsogtgerel. The Lichnerowicz equation on compact manifolds with boundary. Class. Quantum Grav., 30:205011, 2013. Available as arXiv:1306.1801 [gr-qc].
- [9] D. Maxwell. Solutions of the Einstein constraint equations with apparent horizon boundaries. Comm. Math. Phys., 253(3):561-583, 2005.
- [10] D. Maxwell. A class of solutions of the vacuum Einstein constraint equations with freely specified mean curvature. Math. Res. Lett., 16(4):627-645, 2009.
- [11] C. Meier, M. Holst, and G. Tsogtgerel. Non-CMC solutions to the einstein constraint equations with apparent horizon boundaries. 2013. Available as arXiv:1310.2302v1 [gr-qc].
- [12] R. M. Wald. General Relativity. University of Chicago Press, Chicago, IL, 1984.

E-mail address: mholst@math.ucsd.edu

E-mail address: clmeier@math.ucsd.edu

DEPARTMENT OF MATHEMATICS, UNIVERSITY OF CALIFORNIA SAN DIEGO, LA JOLLA CA 92093