OBSTRUCTION-FLAT ASYMMPTOTICALLY
LOCALLY EUCLIDEAN METRICS

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ABSTRACT. We show that any asymptotically locally Euclidean (ALE) metric which is obstruction-flat or extended obstruction-flat must be ALE of a certain optimal order. Moreover, our proof applies to very general elliptic systems and in any dimension \( n \geq 3 \). The proof is based on the technique of Cheeger-Tian for Ricci-flat metrics. We also apply this method to obtain a singularity removal theorem for (extended) obstruction-flat metrics with isolated \( C^0 \)-orbifold singular points.

1. Introduction

We first recall the definition of an ALE metric.

**Definition 1.1.** A complete Riemannian manifold \((M, g)\) is called *asymptotically locally Euclidean* or ALE of order \( \tau \) if it has finitely many ends, and for each end there exists a finite subgroup \( \Gamma \subset SO(n) \) acting freely on \( \mathbb{R}^n \setminus B(0, R) \) and a diffeomorphism \( \Psi : M \setminus \bar{K} \to (\mathbb{R}^n \setminus B(0, R))/\Gamma \) where \( \bar{K} \) is a subset of \( M \) containing all other ends, and such that under this identification,

\[
\begin{align*}
(1.1) & \quad (\Psi_* g)_{ij} = \delta_{ij} + O(r^{-\tau}), \\
(1.2) & \quad \partial^k (\Psi_* g)_{ij} = O(r^{-\tau-k}),
\end{align*}
\]

Date: June 06, 2011. Revised September 28, 2011.
Research partially supported by NSF Grant DMS-0804042.
for any partial derivative of order $k$, as $r \to \infty$, where $r$ is the distance to some fixed basepoint. We say that $(M, g)$ is ALE of order 0 if we can find a coordinate system as above with $(\Psi \ast g)_{ij} = \delta_{ij} + o(1)$, and $\partial^{|k|} (\Psi \ast g)_{ij} = o(r^{-k})$ for any $k \geq 1$ as $r \to \infty$.

ALE spaces are ubiquitous in modern geometric analysis, and we do not attempt to give a complete list of references here. A crucial result in the Ricci-flat case was obtained by Cheeger-Tian: if $(M^n, g)$ is Ricci-flat ALE of order 0, there exists a change of coordinates at infinity so that $(M^n, g)$ is ALE of order $n$, where $n$ is the dimension \[CT94\]. This generalized and improved earlier work of Bando-Kasue-Nakajima \[BKN89\], who employed improved Kato inequalities together with a Moser iteration argument. The Cheeger-Tian method has the advantage of finding the optimal order of curvature decay, without relying on Kato inequalities.

Another interesting class of metrics is that of Bach-flat scalar-flat ALE metrics in dimension 4, or more generally any metric satisfying a system of the form

\[(1.3) \quad \Delta Ric = Rm * Ric,\]

where the right hand side is shorthand for a contraction of the full curvature tensor with the Ricci tensor. In the case of anti-self-dual scalar-flat metrics, or scalar-flat metrics with harmonic curvature, it was proved in \[TV05a\] that such spaces are ALE of order $\tau$ for any $\tau < 2$, using the technique of Kato inequalities. Subsequently, this was generalized to Bach-flat metrics and metrics with harmonic curvature in dimension 4 in \[Str10\], using the Cheeger-Tian technique. In this paper, we will simplify and generalize the Streets argument so that it also works for higher order systems and yields the optimal ALE order. A simplification from \[Str10\] is that we do not need to perform the entire radial separation of variables on symmetric tensors to obtain the optimal decay rate. Rather, we show this optimal decay can be obtained directly in Euclidean coordinates without running into very complicated formulas in radial coordinates, see Proposition \[2.2\].

1.1. The ambient obstruction tensor. Let $(M^n, g)$ be an $n$-dimensional Riemannian manifold, where $n > 2$. Recall that the curvature tensor admits the decomposition

\[(1.4) \quad Rm = W + A_g \otimes g,\]

where $W$ is the Weyl tensor, $\otimes$ is the Kulkarni-Nomizu product, $A_g$ is the Schouten tensor defined as

\[(1.5) \quad A_g = \frac{1}{n-2} \left( Ric - \frac{R}{2(n-1)} g \right),\]

where $R$ denotes the scalar curvature. Define the $n$-dimensional Bach Tensor by (see \[CF08, GH05\])

\[(1.6) \quad B_{ij} = \Delta A_{ij} - \nabla^k \nabla_i A_{kj} + A^{kl} W_{ikjl},\]

where $\Delta$ denotes the rough Laplacian (our convention is to use the analyst’s Laplacian).
If the dimension $n$ is even, then the ambient obstruction tensor introduced in [FG85, FG07], and denoted by $O$ is a symmetric $(0, 2)$-tensor that has the following properties:

1. $O(g)$ is trace-free.
2. If $n = 4$, $O_{ij}(g)$ equals $B_{ij}(g)$ where $B_{ij}(g)$ is the Bach tensor of $g$.
3. If $g$ is conformal to an Einstein metric then $O(g) = 0$.
4. $O(g)$ has an expansion of the form

\[ O_{ij} = \Delta^{\frac{n}{2} - 2} B_{ij} + l.o.t., \]

where $l.o.t.$ denotes quadratic and higher order terms in curvature involving fewer derivatives.

5. $O(g)$ is variational, in fact $O(g)$ is the gradient of the functional

\[ \mathcal{F}(g) = \int_M Q_g dV_g, \]

where $Q_g$ is the $Q$-curvature of $g$. In particular, $O(g)$ is divergence-free.

1.2. Extended obstruction tensors. If $(M^n, g)$ is even-dimensional, there is also a family of symmetric $(0, 2)$-tensors called extended obstruction tensors introduced in [Gra09] and denoted by $\Omega^{(k)}(g)$ where $1 \leq k \leq \frac{n}{2} - 2$ which have the following properties:

1. $\Omega^{(k)}(g)$ is trace-free.
2. When the dimension $n$ is seen as a formal parameter, $\Omega^{(k)}(g)$ has a pole at $n = 2(k + 1)$, and its residue at $n = 2(k + 1)$ is a multiple of the obstruction tensor in that dimension, for example,

\[ \Omega^{(1)} = \frac{1}{4 - n} B_{ij}, \]

and when $n = 4$, $B_{ij}$ equals the obstruction tensor.
3. If $(M, g)$ is locally conformally flat then $\Omega^{(k)}(g) = 0$.
4. $\Omega^{(k)}(g)$ has an expansion of the form

\[ \Omega^{(k)}_{ij} = \frac{1}{(4 - n)(6 - n) \ldots (2k - n - 2)} \Delta^{k - 1} B_{ij} + l.o.t., \]

where $l.o.t.$ denotes quadratic and higher order terms in curvature involving fewer derivatives.

To simplify notation, we define $\Omega^{(k)}(g) = O(g)$ for $k = \frac{n}{2} - 1$.

The main theorem in this paper gives the optimal decay rate for obstruction-flat or extended obstruction-flat scalar-flat ALE metrics:

**Theorem 1.2.** Let $(M^n, g)$ be even-dimensional, scalar-flat, and $\Omega^{(k)}$-flat for some $k$ with $1 \leq k \leq \frac{n}{2} - 1$. If $(M^n, g)$ is ALE of order zero, then there exists a change of coordinates at infinity so that $g$ is ALE of order $n - 2k$. 
The method of Cheeger-Tian is to show that after a suitable change of coordinates, $g$ may be written as $g = g_0 + h$, where $g_0$ is Euclidean, and $h$ is divergence-free. One then considers the linearization of the (extended) obstruction tensor at the flat metric. In the divergence-free gauge, this becomes a power of the Laplacian (the trace is controlled using the scalar-flat condition). An analysis of the decay rates of solutions of the gauged linearized equation, together with an estimate on the nonlinear terms in the equation, then yields Theorem 1.2.

The main technical complication is that the assumption of ALE of order 0 does not directly yield a divergence-free gauge. As in [CT94], we obtain initially a modified divergence-free gauge $\delta_t h = 0$ (see Section 5). In this gauge, we must rule out certain solutions of the linearized equation which we call degenerate solutions (see Definition 5.1). Once these degenerate solutions are ruled out, we are able to find a change of coordinates so that $(M, g)$ is ALE of order $\beta > 0$. This step requires a technique of Leon Simon called the Three Annulus Lemma, which was also employed by Cheeger-Tian [Sim85, CT94]. We generalize this technique so that it applies to higher-order equations. For this step, we show that Turan’s Lemma implies the necessary estimates, which we prove in the Appendix. Once this step is complete, it is relatively easy to find a divergence-free gauge using standard Fredholm Theory, and then to prove the optimal decay order. This work is carried out in Sections 2–6.

Remark 1.3. In the case of the obstruction tensor, which is conformally invariant, one may obtain many examples through the following construction. Let $(M^n, g)$ be an even-dimensional compact Einstein manifold with positive scalar curvature, and let $G_x$ denote the Green’s functions of the conformal Laplacian at a point $x$. The metric $\hat{g} = G_x^p g$, where $p = \frac{4}{n-2}$, is asymptotically flat and scalar-flat [LP87]. Since Einstein spaces are obstruction-flat [GH05, Theorem 2.1], $\hat{g}$ is also obstruction-flat and asymptotically flat of order at least 2. If $(M^n, g)$ is instead locally conformally flat with positive scalar curvature, the same construction yields an $\Omega^{(1)}$-flat asymptotically flat space of order at least $n - 2$.

Our method applies to much more general systems than just the obstruction tensors, and works in any dimension $n \geq 3$. Given two tensor fields $A, B$, the notation $A \ast B$ will mean a linear combination of contractions of $A \otimes B$ yielding a symmetric 2-tensor.

Theorem 1.4. Let $k = 1$ if $n = 3$, or $1 \leq k \leq \frac{n}{2} - 1$ if $n \geq 4$. Assume that $(M, g)$ is scalar-flat, ALE of order 0, and satisfies

$$\Delta_g^k Ric = \sum_{j=2}^{k+1} \alpha_1 + \ldots + \alpha_j = 2(k+1) - 2j \nabla_g^{\alpha_1} Rm \ast \ldots \ast \nabla_g^{\alpha_j} Rm. \tag{1.9}$$

Then $(M, g)$ is ALE of order $n - 2k$.

For $k = 1$, this is simply

$$\Delta Ric = Rm \ast Rm. \tag{1.10}$$
We emphasize that this is more general than (1.3), since the right hand side is allowed to be quadratic in the full curvature tensor. This is satisfied in particular by scalar-flat Kähler metrics and metrics with harmonic curvature in any dimension, and also anti-self-dual metrics in dimension 4. These special cases were previously considered in [Che09] using improved Kato inequalities and a Moser iteration technique. We emphasize that our argument yields the optimal decay rate without requiring any improved Kato inequalities, and therefore applies to the more general system (1.9). The optimal decay for scalar-flat anti-self-dual ALE metrics was previously considered in [CLW08, Proposition 13]. The case of extremal Kähler ALE metrics was considered in [CW11]. As mentioned above, the cases of Bach-flat metrics and metrics with harmonic curvature in dimension 4 were considered in [Str10]. However, we note that (1.10) is more general than (1.3).

Remark 1.5. We do not need such a strong requirement on the decay of partial derivatives of arbitrarily high order in (1.2) in Definition 1.1, but have assumed this here in the introduction for simplicity of stating the result. We only need to assume this up to a finite number of partial derivatives, see Remark 5.9.

1.3. Singularity removal. The methods used to prove the above results can also be applied to analyze isolated singularities. Similar results were proved in [BKN89, Che09, CW11, CLW08, Str10, Tia90, TV05b]. We next recall the definition of a $C^0$-orbifold point.

Definition 1.6. Let $g$ be a metric defined on $B_\rho(0) \setminus \{0\}$, where $B_\rho(0)$ is a metric ball in a flat cone. We say that the origin is a $C^0$-orbifold point if there exists a coordinate system around the origin such that

$$g_{ij} = \delta_{ij} + o(1),$$

$$\partial^l g_{ij} = o(r^{-|l|}),$$

for any multi-index $l$ with $|l| \geq 1$ as $r \to 0$. We say that the origin is a smooth orbifold point, if after lifting to the universal cover of $B_\rho(0) \setminus \{0\}$, the metric extends to a smooth metric over the origin, after diffeomorphism.

Remark 1.7. As in the ALE case, we will not need to assume (1.12) for partial derivatives of arbitrarily high order, only up to a certain finite number of derivatives, see Remark 7.2.

Applying the Cheeger-Tian technique directly to the singularity, we obtain the following.

Theorem 1.8. Let $B_\rho(0)$ be as above and even-dimensional, and let $g$ be (extended) obstruction-flat in $B_\rho(0) \setminus \{0\}$ with constant scalar curvature. If the origin is a $C^0$-orbifold point for $g$, then the metric extends to a smooth orbifold metric in $B_\rho(0)$.

As in the ALE case, this theorem also applies to much more general higher-order systems:
Theorem 1.9. Let $k = 1$ if $n = 3$, or $1 \leq k \leq \frac{n}{2} - 1$ if $n \geq 4$. Assume that $(B_\rho(0) \setminus \{0\}, g)$ has constant scalar curvature and satisfies
\begin{equation}
\Delta^k_g \text{Ric} = \sum_{j=2}^{k+1} \sum_{\alpha_1 + \ldots + \alpha_j = 2(k+1)-2j} \nabla^\alpha \text{Rm} \ast \ldots \ast \nabla^\alpha \text{Rm}.
\end{equation}

If the origin is a $C^0$-orbifold point for $g$, then the metric extends to a smooth orbifold metric in $B_\rho(0)$.

Theorem 1.9 will be proved in Section 7, the proof of which uses the same method as that of Theorem 1.2, with a few minor modifications.

1.4. Acknowledgements. The authors would like to thank Michael Singer and Gang Tian for invaluable assistance from the inception of this project. We would also like to thank Robin Graham for crucial help with the structure of the obstruction tensors, Jeff Streets for numerous helpful discussions, and Sigurd Angenent and Mikhail Feldman for several helpful comments regarding elliptic theory. We are greatly indebted to Fedja Nazarov for explaining how Turan’s Lemma could be used to prove Lemma 5.2. Finally, we would like to thank the anonymous referee for numerous suggestions which greatly improved the exposition of the paper.

2. Linearized obstruction tensor

We begin with some notation: in the following $\delta$ will denote the divergence operator, which can act either on a symmetric 2-tensor $h$, or on a 1-form $\omega$. In the former case $\delta h = \nabla^i h_{ij}$ and in the latter case $\delta \omega = \nabla^i \omega_i$. The $L^2$-adjoint of $\delta$ will be denoted by $\delta^*$, which is $-(1/2)\mathcal{L}$, where $\mathcal{L}$ is the Lie derivative operator, defined by $(\mathcal{L} \omega)_{ij} = \nabla_i \omega_j + \nabla_j \omega_i$. The trace of a symmetric 2-tensor $h$ will be denoted by $\text{tr}(h)$.

We now analyze the linearizations of (1.7) and (1.8). Note that from the dependence of the lower order terms in both equations on the curvature tensor it follows that the linearized equations $(\Omega^{(k)})'_{g_0}(h) = 0$ for $1 \leq k \leq \frac{n}{2} - 1$ at a flat metric $g_0$ are equivalent to $\Delta^{k-1} B'_{g_0}(h) = 0$. With this observation we have the following, which holds in any dimension $n \geq 3$:

Proposition 2.1. At a flat metric $g_0$ we have
\begin{equation}
\Delta^{k-1} \left( B'_{g_0}(h) \right) = \Delta^{k-1} \left( -\frac{1}{2(n-2)} \Delta^2 h - \frac{1}{2(n-1)(n-2)} \nabla^2 \text{tr}(h) - \frac{1}{2(n-1)} \nabla^2 \delta h - \frac{1}{(n-2)} \Delta^* \delta h + \frac{1}{2(n-1)(n-2)} \left( \Delta^2 \text{tr}(h) - \Delta \delta h \right)_{g_0} \right).
\end{equation}

Proof. The Bach tensor can be written as $B_{ij} = \Delta A_{ij} - \nabla^k \nabla_i A_{kj} + A^{kl} W_{ikjl}$, and at a flat metric $(A^{kl} W_{ikjl})'_{g_0}(h) = 0$ for any $h \in S^2(T^*\mathbb{R}^n)$, so then
\begin{equation}
B'_{ij}(h) = \Delta A'_{ij}(h) - \nabla^k \nabla_i A'_{jk}(h).
\end{equation}
By the Bianchi identity and using again that \(g_0\) is flat we have
\[
-\nabla^k \nabla_i (A'_{g_0})_{jk}(h) = -\nabla_i \nabla^k (A_{g_0})'_{jk}(h) = \frac{1}{2(n-1)} \nabla_i \nabla_j R'_{g_0}(h)
\]
(2.2)
\[
= \frac{1}{2(n-1)} \nabla_i \nabla_j (\Delta \text{tr}(h) - \delta(\delta h)).
\]

On the other hand
\[
\Delta A'_{g_0}(h) = \frac{1}{(n-2)} \Delta \left( \text{Ric}'_{g_0}(h) - \frac{R'_{g_0}}{2(n-1)}(h)g_0 \right)
\]
(2.3)
\[
= \frac{1}{(n-2)} \Delta \left( -\frac{1}{2} \Delta h - \frac{1}{2} \nabla^2 \text{tr}(h) - \delta^* \delta h + \frac{1}{2(n-1)} [\Delta \text{tr}(h) - \delta(\delta h)] g_0 \right).
\]

Adding (2.2) and (2.3) together and taking \(\Delta^{k-1}\) we obtain (2.1).

The linearized equations are not strictly elliptic due to the diffeomorphism invariance of the (extended) obstruction-flat equations. Note, however, that if \(h\) is divergence free and if the trace of \(h\) is harmonic then from (2.1) the linearized equations reduce to
\[
\Delta^{k+1} h = 0,
\]
(2.4)
\[
\Delta \text{tr}(h) = 0,
\]
(2.5)
\[
\delta h = 0,
\]
(2.6)
and (2.4) is an elliptic equation on \(h\). We point out that we will not be able to prescribe that \(h\) be divergence-free at first, so we will follow the approach in [CT94] and introduce a modified divergence operator in Section 3.1. The proof of Theorem 1.2 relies on the following crucial proposition.

**Proposition 2.2.** Let \(n > 2\), and \(h \in S^2(T^*\mathbb{R}^n)\) be a solution on \(\mathbb{R}^n \setminus B_\rho(0)\) for some \(\rho > 0\) of the system
\[
\Delta^{k+1} h = 0,
\]
(2.7)
\[
\delta h = 0,
\]
(2.8)
for \(1 \leq k \leq \frac{n}{2} - 1\) (or for \(k = 1\) when \(n = 3\)) satisfying \(h = O(|x|^{1-\epsilon})\) with \(\epsilon > 0\). Then, \(h\) can be written as
\[
h = h_c + O \left( |x|^{-n+2k} \right),
\]
(2.9)
where \(h_c\) is constant. The result also holds for \(k = 0\) if in addition we assume that \(\text{tr}(h) = 0\).

**Proof.** Consider first the case \(1 \leq k \leq \frac{n}{2} - 2\). Since the components of \(h\) satisfy the scalar equation \(\Delta^{k+1} f = 0\), using a classical expansion, we may write a solution \(h\) of
(2.7)-(2.8) as

\[ h(x) = h_c + \sum_{l=0}^{\infty} h_l(x), \]

where \( h_c \) is constant and each \( h_l \) is a homogeneous solution of (2.7)-(2.8) of degree \( 2(k+1) - n - l \). If \( h_l \) is one of such solutions, we can write

\[ h_l(x) = |x|^{2(k+1)-n-l} \sum_{(s,j) \in S_{k,l}} h_{s,j}(x), \]

where \( S_{k,l} = \{(s,j) : 2(j+1) - s = 2(k+1) - l, s \geq 0, 0 \leq j \leq k\} \), and the components of \( h_{s,j}(x) \) in (2.11) are spherical harmonics of degree \( s \). Note that if \((s,j) \in S_{k,l}\), then \( \Delta^{j+1}(|x|^{2(k+1)-n-l}h_{s,j}(x)) = 0 \). In order to prove the claim for \( 1 \leq k \leq \frac{n}{2} - 2 \), it suffices to show that \( h_0 = h_1 \equiv 0 \). For \( l = 0 \), we have

\[ (h_0)_{ij}(x) = |x|^{2(k+1)-n}c_{ij}, \]

where \( c_{ij} \) is constant. From (2.8) we obtain

\[ \sum_{i=1}^{n} c_{ij}x_i = 0, \]

for each \( j \) and this clearly implies that \( h_0 \) is identically zero. For \( h_1 \) we have

\[ (h_1)_{ij}(x) = |x|^{2(k+1)-n}u_{ij}\left(\frac{x}{|x|^2}\right), \]

where \( u_{ij}(x) \) is a homogeneous polynomial of degree 1 that we will write as

\[ u_{ij}(x) = \sum_{i=1}^{n} A_{ij}x_i. \]

From (2.8) we have for every \( j \)

\[ 0 = \sum_{i=1}^{n} \partial_i \left(|x|^{2(k+1)-n}u_{ij}(x/|x|^2)\right) = -\sum_{i=1}^{n} (n - 2k)x_i|x|^{2(k-1)-n}u_{ij}(x) + |x|^{2k-n}\sum_{i=1}^{n} \partial_i u_{ij}(x), \]

which becomes

\[ (n - 2k) \sum_{i=1}^{n} \sum_{l=1}^{n} A_{ij}x_i x_l = |x|^2 \sum_{i=1}^{n} A_{ij}. \]

For fixed \( j \) and every \( x \in \mathbb{R}^n \), with \( x \neq 0 \). If in (2.13) we let \( x \) be the vector with coordinates \( x_i = \delta_{ip} \) for fixed \( p \), one obtains the identity

\[ A_{pjp} = \frac{1}{n - 2k} \sum_{l=1}^{n} A_{ijl} \text{ for fixed } j. \]
An obvious consequence of (2.14) is that $A_{pjp}$ is independent of $p$ for fixed $j$, in particular, for every $p$ and $j$

\begin{equation}
A_{pjp} = \frac{n}{n-2k} A_{pjp},
\end{equation}

and then $A_{pjp} = 0$ since $k \neq 0$. For the components of the form $A_{ijm}$ with $l \neq m$, the coefficient of $x_l x_m$ in the left-hand side of (2.13) is $(n-2k)(A_{ijm} + A_{mjl})$ while in the right-hand side there are no off-diagonal terms, so we conclude that

\begin{equation}
A_{ijm} = -A_{mjl} \text{ for } l \neq m,
\end{equation}

If $l, j, m$ are all different we obtain from the symmetry of $A_{ijm}$ in $l, j$ and from (2.16) the identity

$$A_{ijm} = -A_{mjl} = -A_{ijl} = A_{mlj} = -A_{jlm} = -A_{ijm},$$

therefore, in this case $A_{ijm} = 0$. For the components of the form $A_{llj}$ and $A_{jlj}$ when $l \neq j$, it is easy to see that $A_{llj} = -A_{jlj}$ and $A_{jlj} = A_{llj}$ and as we saw above this implies that both components are zero, so we conclude that all polynomials $u_{ij}$ are identically zero.

For the case $k = \frac{n}{2} - 1$ there is only one difference with the argument above: $h_0(x)$ is logarithmic, i.e., a solution of the form $h_{ij}(x) = \log(|x|) c_{ij}$ with $c_{ij}$ constant, however the condition $\delta h = 0$ implies that $\sum_{j=1}^{n} c_{ij} x_j = 0$ for every $i$ and hence $c_{ij} = 0$ for all $i, j$ so this solution in fact does not occur.

For the case $k = 1$ and $n = 3$ we write $h$ as

\begin{equation}
h(x) = \sum_{l=0}^{\infty} h_l(x),
\end{equation}

where $h_l(x)$ is a homogeneous solution of degree $-l$ of $\Delta^2 h(x) = 0$ on $\mathbb{R}^n \setminus \{0\}$. In this case, the solution $h_0(x)$ has the form $h_0(x) = h_C + h_{0,1}(x)$ where the components of $h_C$ are constant and the components of $h_{0,1}$ are spherical harmonics of degree 1. The solutions $h_l(x)$ with $l \geq 1$ have the form $h_l(x) = |x|^{-l} (h_{l,t-1}(x) + h_{l,t+1}(x))$ where the components of $h_{l,t\pm1}(x)$ are spherical harmonics of degree $l \pm 1$. We only have to prove that if $\delta h_{0,1}(x) = 0$ on $\mathbb{R}^n \setminus \{0\}$ then $h_{0,1}(x) \equiv 0$. For that purpose write the components of $h_{0,1}(x)$ as $(h_{0,1})_{ij}(x) = u_{ij} \left( \frac{x}{|x|} \right)$ where $u_{ij}(x)$ are linear functions given by (2.12). The condition $\delta h_{0,1}(x) = 0$ for all $x \neq 0$ becomes

$$\sum_{i=1}^{3} \sum_{l=1}^{3} A_{ijl} x_i x_l = |x|^2 \sum_{i=1}^{3} A_{iji},$$

and we can argue as in the case $1 \leq k \leq \frac{n}{2} - 1$ to conclude that $A_{ijl} = 0$ for all $i, j, l = 1, 2, 3$.

The above proof can be extended to the case of $k = 0$ provided the trace vanishes. However, we omit the proof, and instead refer the reader to the proof given in [C194] page 538, which is an alternative argument using Obata’s Theorem. □
3. **Nonlinear terms in the obstruction-flat systems**

In this section we derive an expression for the error terms in the linearization of the (extended) obstruction tensors, i.e., the difference

\[ \Omega^{(k)}(g_0 + h) - \Omega^{(k)}(g_0) - \left( \Omega^{(k)} \right)'_{g_0}(h), \]

where \( g_0 \) is a flat metric in \( \mathbb{R}^n \). Given two tensor fields \( A, B \) by \( A \ast B \) we mean a linear combination of contractions of \( A \otimes B \) using the metric \( g_0 \), and for a positive integer \( j \), \( A^{-j} \ast B \) means contractions of \( j \) copies of the inverse of \( A \) with \( B \).

**Proposition 3.1.** Let \( g_0 \) be a flat metric on \( \mathbb{R}^n \) and let \( h \in S^2(T^*\mathbb{R}^n) \) be such that \( g_0 + h \) is another Riemannian metric on \( \mathbb{R}^n \). For the (extended) obstruction tensors \( \Omega^{(k)} \) with \( 1 \leq k \leq \frac{n}{2} - 1 \), we have

\[ \Omega^{(k)}(g_0 + h) = \left( \Omega^{(k)} \right)'_{g_0}(h) - (g_0 + h)^{-1} \ast h \ast \nabla^{2(k+1)}_{g_0} h \]

\[ + \sum_{j=2}^{\mathcal{I}_k} (g_0 + h)^{-j} \ast \left( \sum_{\alpha_1 + \cdots + \alpha_j = 2(k+1)} \nabla^{\alpha_1}_{g_0} h \ast \cdots \ast \nabla^{\alpha_j}_{g_0} h \right), \]

for some integer \( \mathcal{I}_k > 2(k+1) \). For the scalar curvature we have,

\[ R(g_0 + h) = R'_{g_0}(h) + (g_0 + h)^{-1} \ast h \ast \nabla^2_{g_0} h + (g_0 + h)^{-2} \ast (h \ast \nabla^2_{g_0} h + \nabla_{g_0} h \ast \nabla_{g_0} h) \]

\[ + (g_0 + h)^{-3} \ast (\nabla_{g_0} h \ast \nabla_{g_0} h \ast h). \]

**Proof.** For any tensor \( T \), we have

\[ \nabla_{g_0 + h} T = \nabla_{g_0} T + (g_0 + h)^{-1} \ast \nabla_{g_0} h \ast T. \]

From this it follows that for the (1,3) curvature tensor

\[ Rm(g_0 + h) = (g_0 + h)^{-1} \ast \nabla^2_{g_0} h + (g_0 + h)^{-2} \ast \nabla_{g_0} h \ast \nabla_{g_0} h. \]

It follows that the \( Ric(g_0 + h) \) has an expansion similar to (3.3). For the scalar curvature, on the other hand, we have

\[ R(g_0 + h) = (g_0 + h)^{-1} \ast (Ric(g_0 + h)) \]

\[ = (g_0 + h)^{-1} \ast ((g_0 + h)^{-1} \ast \nabla^2_{g_0} h + (g_0 + h)^{-2} \ast \nabla_{g_0} h \ast \nabla_{g_0} h). \]

Using the identity

\[ (g_0 + h)^{-1} - g_0^{-1} = -g_0^{-1} \ast h \ast (g_0 + h)^{-1}, \]

we obtain (we will omit writing \( g_0^{-1} \) from now on)

\[ R(g_0 + h) = (g_0 + h)^{-1} \ast \nabla^2_{g_0} h + (g_0 + h)^{-2} \ast (h \ast \nabla^2_{g_0} h + \nabla_{g_0} h \ast \nabla_{g_0} h) \]

\[ + (g_0 + h)^{-3} \ast h \ast \nabla_{g_0} h \ast \nabla_{g_0} h, \]
and another application of (3.7) to the leading term yields

\begin{equation}
R(g_0 + h) = \nabla^2_{g_0} h + (g_0 + h)^{-1} \ast h \ast \nabla^2_{g_0} h + (g_0 + h)^{-2} \ast (h \ast \nabla^2_{g_0} h + \nabla_{g_0} h \ast \nabla_{g_0} h) \\
+ (g_0 + h)^{-3} \ast h \ast \nabla_{g_0} h \ast \nabla_{g_0} h.
\end{equation}

(3.8)

Since the only term in (3.8) that contributes to the linearization of \( R \) at \( g_0 \) is \( \nabla^2_{g_0} h \), equation (3.3) follows. In order to find a similar expansion for the Bach tensor \( B(g_0 + h) \) we note that for any metric \( g \), \( B(g) \) can be written schematically as

\begin{equation}
B(g) = \Delta_g \text{Ric}_g + \nabla^2_g R_g + (\Delta_g R_g) g + Rm_g \ast Rm_g.
\end{equation}

(3.9)

We first consider the term \( \Delta_g \text{Ric}_g \) at \( g = g_0 + h \) in (3.9). Taking a covariant derivative, again assuming that \( g_0 \) is flat, we have

\begin{equation}
\nabla_{g_0 + h} \text{Ric}(g_0 + h) = (g_0 + h)^{-1} \ast (\nabla^3_{g_0} h) + (g_0 + h)^{-2} \ast \nabla^2_{g_0} h \ast \nabla_{g_0} h \\
+ (g_0 + h)^{-3} \ast \nabla_{g_0} h \ast \nabla_{g_0} h \ast \nabla_{g_0} h.
\end{equation}

(3.10)

Taking another covariant derivative and one contraction with respect to \( g_0 + h \), we obtain

\begin{equation}
\Delta_{g_0 + h} \text{Ric}(g_0 + h) = (g_0 + h)^{-1} \ast (\nabla^4_{g_0} h) \\
+ (g_0 + h)^{-2} \ast (\nabla^2_{g_0} h \ast \nabla^2_{g_0} h + \nabla^3_{g_0} h \ast \nabla_{g_0} h) \\
+ (g_0 + h)^{-3} \ast \nabla_{g_0} h \ast \nabla_{g_0} h \ast \nabla_{g_0} h \\
+ (g_0 + h)^{-4} \ast \nabla_{g_0} h \ast \nabla_{g_0} h \ast \nabla_{g_0} h \\
+ (g_0 + h)^{-5} \ast \nabla_{g_0} h \ast \nabla_{g_0} h \ast \nabla_{g_0} h \ast \nabla_{g_0} h.
\end{equation}

(3.11)

Using (3.7) twice (as we did above for \( R \)), we conclude that

\begin{equation}
\Delta_{g_0 + h} \text{Ric}(g_0 + h) = \Delta_{g_0} \text{Ric}'_{g_0}(h) + (g_0 + h)^{-1} \ast h \ast \nabla^4_{g_0} h \\
+ \sum_{j=2}^{5} \sum_{\alpha_1 + \ldots + \alpha_j = 4} (g_0 + h)^{-j} \ast \nabla^{\alpha_1}_{g_0} h \ast \ldots \ast \nabla^{\alpha_j}_{g_0} h.
\end{equation}

(3.12)

From similar computations for the other terms in (3.9) we conclude that

\begin{equation}
B(g_0 + h) = B'_{g_0}(h) + (g_0 + h)^{-1} \ast h \ast \nabla^4_{g_0} h \\
+ \sum_{j=2}^{6} \sum_{\alpha_1 + \ldots + \alpha_j = 4} (g_0 + h)^{-j} \ast \nabla^{\alpha_1}_{g_0} h \ast \ldots \ast \nabla^{\alpha_j}_{g_0} h.
\end{equation}

(3.13)

For any \( k \geq 1 \), a similar argument shows that

\begin{equation}
\Delta_{g_0 + h}^k B(g_0 + h) = \Delta_{g_0}^k B'_{g_0}(h) + (g_0 + h)^{-1} \ast h \ast \nabla^{2(k+1)}_{g_0} h \\
+ \sum_{j=2}^{2k} (g_0 + h)^{-j} \ast \left( \sum_{\alpha_1 + \ldots + \alpha_j = 2(k+1)} \nabla^{\alpha_1}_{g_0} h \ast \ldots \ast \nabla^{\alpha_j}_{g_0} h \right).
\end{equation}

(3.14)

where \( I_k \) is some integer with \( I_k > 2(k+1) \).
Next, using scaling arguments it is clear that the terms \(l.o.t.\) in (1.7) and (1.8) have the form
\[
(3.15) \quad \sum_{j=2}^{k+1} \sum_{\alpha_1 + \ldots + \alpha_j = 2(k+1)-2j} \nabla^{\alpha_j} Rm \ast \ldots \ast \nabla^{\alpha_1} Rm,
\]
see for example the proof of Theorem 2.1 in [GH05]. Equation (3.2) follows from combining (1.7) or (1.8) in Subsection 1.2 with (3.14), since the terms in (3.15) admit a similar expansion as the nonlinear terms in (3.14) (and they do not affect the linearization). □

Next, defining
\[
(3.16) \quad c_{n,k} = \begin{cases} 
\frac{1}{(4-n)(6-n)\ldots(2k-n)} & \text{if } 1 \leq k \leq \frac{n}{2} - 2 \\
1 & \text{if } k = \frac{n}{2} - 1,
\end{cases}
\]
from Proposition 3.1, we may write
\[
(3.17) \quad \Omega^{(k)}(g_0 + h) = c_{n,k} \Delta^{k-1} B'_{g_0}(h) + F^{(k)}(h, g_0),
\]
where \(F^{(k)}(h, g_0)\) is the remainder in (3.2). For the scalar curvature we will write
\[
(3.18) \quad R(g_0 + h) = R'_{g_0}(h) + F'(h, g_0),
\]
where \(F'(h, g_0)\) is the error term in (3.3). From now on, we will use \(\nabla\) to denote \(\nabla_{g_0}\), therefore all operators \(\Delta, \delta, \text{tr}\) are taken with respect to \(g_0\).

3.1. Scalar-flat condition and the modified equation. In order to address the difficulty of not initially being able to prescribe \(h\) to be divergence-free, we follow [CT94] and introduce a modified divergence operator given by
\[
(3.19) \quad \delta_t h = \delta h - t i_{r-1} \frac{\partial}{\partial r} h,
\]
and we will show in Section 5 that we can find a gauge where \(\delta_t h = 0\). A difference with the approach in [CT94] is that the obstruction-flat systems are not elliptic even if we are able to prescribe \(\delta_t h = 0\) because \(\Delta^{k-1} B'_{g_0}(h)\) is traceless regardless of the gauge condition. Note that if \(\delta_t h = 0\) we obtain
\[
(3.20) \quad \Delta^{k-1} B'_{g_0}(h) = \Delta^{k-1} \left( -\frac{1}{2(n-2)} \Delta^2 h - \frac{1}{2(n-1)(n-2)} \nabla^2 \Delta \text{tr}(h) \\
- \frac{t}{2(n-1)} \nabla^2 \delta i_{r-1} \frac{\partial}{\partial r} h - \frac{t}{(n-2)} \Delta \delta^* i_{r-1} \frac{\partial}{\partial r} h \\
+ \frac{1}{2(n-2)(n-1)} \left( \Delta^2 \text{tr}(h) - t \Delta \delta \left( i_{r-1} \frac{\partial}{\partial r} h \right) \right) g_0 \right).
\]
At this point we use the scalar-flat condition on \(g_0 + h\). Assuming again that \(\delta_t h = 0\), the linearization of the scalar curvature at \(g_0\) becomes
\[
(3.21) \quad R'_{g_0}(h) = -\Delta \text{tr}(h) + \delta \delta h = -\Delta \text{tr}(h) + t \delta i_{r-1} \frac{\partial}{\partial r} h,
\]
so from the scalar-flat equation we have
\[ \Delta \text{tr}(h) = t \delta i_{r-1, \frac{\partial}{\partial r}} h + F'(h, g_0), \]
where \( F'(h, g_0) \) is the remainder in (3.18). Inserting (3.22) into (3.20) we obtain
\[ \Delta^{k-1} B'_{g_0}(h) = \Delta^{k-1} \left( -\frac{1}{2(n-2)} \Delta^2 h - \frac{t}{2(n-2)} \nabla^2 \delta i_{r-1, \frac{\partial}{\partial r}} h \right) \]
\[ - \frac{t}{(n-2)} \Delta^* i_{r-1, \frac{\partial}{\partial r}} h + \mathcal{E}^{(k)}(h, g_0), \]
where \( \mathcal{E}^{(k)}(h, g_0) \) is given by
\[ \mathcal{E}^{(k)}(h, g_0) = \frac{1}{2(n-1)(n-2)} \left( - \Delta^{k-1} \nabla^2 F'(h, g_0) + \Delta^k F'(h, g_0) g_0 \right). \]

We now define a linear operator \( \mathcal{P}^{(k)}_t \) by
\[ \mathcal{P}^{(k)}_t(h) = \frac{c_{n,k}}{n-2} \Delta^{k-1} \left( -\frac{1}{2} \Delta^2 h - \frac{t}{2} \nabla^2 \delta i_{r-1, \frac{\partial}{\partial r}} h - t \Delta^* i_{r-1, \frac{\partial}{\partial r}} h \right). \]
Clearly, the operator \( \mathcal{P}^{(k)}_t \) is strictly elliptic. From (3.2) and (3.23), if \( \delta_i h = 0 \) and \( g_0 + h \) is scalar-flat, the (extended) obstruction-flat system may be written as
\[ 0 = \mathcal{P}^{(k)}_t h + \mathcal{R}^{(k)}(h, g_0) \text{ for } 1 \leq k \leq \frac{n}{2} - 1, \]
where
\[ \mathcal{R}^{(k)}(h, g_0) = c_{n,k} \mathcal{E}^{(k)}(h, g_0) + F^{(k)}(h, g_0). \]
Writing \( \mathcal{E}^{(k)}(h, g_0) \) schematically as \( \nabla^{2k} F'(h, g_0) \), we easily see using the proof of Proposition 3.1 that \( \mathcal{R}^{(k)}(h, g_0) \) has the same form as the remainder in (3.2).
\[ \mathcal{R}^{(k)}(h, g_0) = (g_0 + h)^{-1} * h * \nabla^{2(k+1)} g_0 h \]
\[ + \sum_{j=2}^{k+1} (g_0 + h)^{-j} * \left( \sum_{\alpha_1 + \ldots + \alpha_j = 2(k+1)} \nabla^{\alpha_1} g_0 h * \ldots * \nabla^{\alpha_j} g_0 h \right), \]
so that none of the terms in \( \mathcal{R}^{(k)}(h, g_0) \) contribute to the linearization. This shows that (3.26) defines a family of elliptic equations on \( \mathbb{R}^n \setminus \{0\} \). This same argument applies equally to the system (1.9). We have proved

**Corollary 3.2.** Let \( 1 \leq k \leq \frac{n}{2} - 1 \). If \( g = g_0 + h \) is scalar-flat and \( \Omega^{(k)} \)-flat with \( \delta_i h = 0 \), then \( h \) satisfies
\[ \mathcal{P}^{(k)}_t h + \mathcal{R}^{(k)}(h, g_0) = 0, \]
where \( \mathcal{P}^{(k)}_t \) is the linear operator given by (3.22) and \( \mathcal{R}^{(k)}(h, g_0) \) is the error term given by (3.27). The operator \( \mathcal{P}^{(k)}_t \) is strictly elliptic.

Let \( k = 1 \) if \( n = 3 \), or \( 1 \leq k \leq \frac{n}{2} - 1 \) if \( n \geq 4 \). If \( g = g_0 + h \) is scalar-flat and solves (1.9) with \( \delta_i h = 0 \), then \( h \) also satisfies (3.28) with a remainder term of the same form.
4. Weighted Hölder and Sobolev Spaces

In this section we introduce weighted spaces that will be useful in the analysis needed to construct divergence-free gauges. We start by reviewing some of the notation in [CT94]. If we write \( \mathbb{R}^n \) as \( \mathbb{R}^n = C(S^{n-1}) \), we define maps \( \psi_a : C(S^{n-1}) \to C(S^{n-1}) \) for \( a > 0 \) given by \( \psi_a(r, x) = (ar, x) \). The weighted Hölder norms are defined as follows: if \( u > 0 \), let \( A_u \) be the natural scaling on tensors of type \((p, q)\), i.e.

\[
A_u = (\psi_{u^{-1}})_* \otimes \cdots \otimes (\psi_{u^{-1}})_* \otimes \psi_u^* \otimes \cdots \otimes \psi_u^* .
\]

Given a tensor \( T \) of type \((p, q)\) we have

\[
|T_{(u,x)}|_{m,\alpha;0} = |u^{p-q}A_u T_{(1,x)}|_{m,\alpha;0},
\]

where \( | \cdot |_{m,\alpha;0} \) is the \( C^{m,\alpha} \)-norm with respect to the flat metric \( g \) at the point \((1, x)\). We can now define Hölder norms by

\[
|T|_{m,\alpha;0} = \sup_{(u,x)} |T_{(u,x)}|_{m,\alpha;0},
\]

and also weighted Hölder norms given by

\[
|T|_{m,\alpha;l} = |r^{-l}T|_{m,\alpha;0},
\]

for any \( l \in \mathbb{R} \). We say that a tensor \( T \) of type \((p, q)\) is in \( \mathcal{T}_{m,\alpha;l}^{p,q} \) if \( |T|_{m,\alpha;l} < \infty \). We use \( A_{c,d}(0) \) to denote the annulus \( A_{c,d}(0) = \{(r, x) | c < r < d\} \). For a tensor \( h \) of type \((u, v)\) and \( \delta \in \mathbb{R} \) we define a weighted \( L^p \) norm by

\[
\|h\|_{L^p_{\delta}} = \left( \int_{\mathbb{R}^n} |h|^p |x|^{-\delta p-n} dx \right)^{\frac{1}{p}},
\]

where \( | \cdot | \) is the usual pointwise norm on tensors of type \((u, v)\). For any nonnegative integer \( m \), we also define weighted Sobolev norms by

\[
\|h\|_{W_{\delta}^m,p,u,v} = \sum_{j=0}^m \|\nabla^j h\|_{L^p_{\delta-j}},
\]

and then \( W_{\delta}^m,p,u,v \) is the space

\[
W_{\delta}^m,p,u,v = \{ h : \|h\|_{W_{\delta}^m,p,u,v} < \infty \}.
\]

For further properties of these weighted Sobolev spaces see [Bar86].

For notational convenience, if \( h \) is a symmetric \((0, 2)\)-tensor we will use \( |||h|||_{a,b} \) to denote the norm \( \|h\|_{L^2_{\delta}(A_{a,b}(0))} \), i.e., the \( L^2 \) norm of \( h \) with weight 0 on the annulus \( A_{a,b}(0) \). Another way to construct the norm \( ||| \cdot ||| \) is as follows: consider the weighted inner product on the slices \((r, S^{n-1})\) given by

\[
\langle\langle h_1, h_2 \rangle\rangle = r^{-(n-1)} \int_{(r, S^{n-1})} \langle h_1, h_2 \rangle dV_{g_{S^{n-1}(r)}},
\]

where \( g_{S^{n-1}(r)} \) is the metric on the sphere of radius \( r \) in \( S^{n-1} \).
where $\langle \cdot, \cdot \rangle$ is the usual pointwise inner product. It follows that
\begin{equation}
\|||h|||^2_{a,b} = \int_a^b r^{-1} ||h||^2 dr,
\end{equation}
where $|| \cdot ||$ is the norm defined by the inner product $\langle \langle \cdot, \cdot \rangle \rangle$ in (4.8). The norms $||| \cdot |||$ are scale invariant in the sense that if $h$ is a $(0, 2)$ tensor and if we let $q = a^{-2} \psi^a h$ then
\begin{equation}
|||h|||_{a,aL} = |||q|||_{1,L}.
\end{equation}
Finally, we say that a $(p, q)$-tensor $T$ is radially parallel if
\begin{equation}
\nabla_{\frac{\partial}{\partial r}} T = 0.
\end{equation}

4.1. **Divergence and the Lie derivative operator.** We now consider the operator $\Box : \Lambda^1(\mathbb{R}^n) \to \Lambda^1(\mathbb{R}^n)$ defined by $\Box \xi = \delta L_\xi g_0$, where $L_\xi$ is the Lie derivative operator. This operator is formally self-adjoint and elliptic. From now on we use $\tilde{\Delta}_H$, $\tilde{\delta}$ and $\tilde{d}$ to denote the Hodge laplacian, trace, divergence and exterior differentiation in the cross section metric $g_{S^{n-1}}$, respectively. Following [CT94, Section 2], if we write $\xi$ in polar coordinates as
\begin{equation}
\xi = f dr + \omega,
\end{equation}
we have
\begin{equation}
L_\xi g_0 = 2f' dr \otimes dr + \left( \omega' - 2r^{-1} \omega + \tilde{d} f \right) \Box dr + L_\omega g_{S^{n-1}} + 2rf g_{S^{n-1}},
\end{equation}
here we use primes to denote differentiation respect to $r$. Also, given a 1-form $\alpha$ we denote by $\alpha \Box dr$ the symmetric product $\alpha \otimes dr + dr \otimes \alpha$. If now $h$ is a symmetric $(0, 2)$-tensor written in polar coordinates as
\begin{equation}
h = h_{00} dr \otimes dr + \alpha \Box dr + B,
\end{equation}
where $B$ is a symmetric $(0, 2)$ tensor whose radial components are zero, the divergence of $h$ is given by
\begin{equation}
\delta h = \left( h'_{00} + (n-1)r^{-1}h_{00} + r^{-2}\tilde{\delta} \alpha - r^{-3} \tilde{\delta} \Box (B) \right) dr
+ \alpha' + (n-1)r^{-1} \alpha + r^{-2} \tilde{\delta} B.
\end{equation}
Combining (4.13) and (4.15), the operator $\Box$ takes the form
\begin{equation}
\Box \xi = \left( 2f'' + 2(n-1)r^{-1}f' + r^{-2} \left( -\tilde{\Delta}_H f - 2(n-1)f \right) + r^{-2} \tilde{\delta} \omega' - 4r^{-2} \tilde{\delta} \omega \right) dr
+ \omega'' + (n-3)r^{-1} \omega' + r^{-2} \left( -\tilde{\Delta}_H \omega + \tilde{d} \omega \right) + \tilde{d} f' + r^{-1}(n+1) \tilde{d} f.
\end{equation}
As pointed out in [CT94, Section 2], any 1-form defined on $\mathbb{R}^n \{0\}$ can be written as an infinite sum of forms of two types
1. Type I: $p(r)\psi$ where $\psi \in \Lambda^1(T^*S^{n-1})$, $\tilde{\delta} \psi = 0$ and $\tilde{\Delta}_H \psi = \mu \psi$,
2. Type II: $r^{-1}l(r)\phi dr + u(r)rd\phi$ where $\phi \in \Lambda^0(T^*S^{n-1})$ and $\tilde{\Delta}_H \phi = \nu \phi$.  


Moreover, the operator $\Box$ preserves these two types of forms. If $\xi$ is a 1-form of type I or II, the equation $\Box \xi = 0$ reduces to solving a second order linear system of ordinary differential equations of at most two equations. In order to see what these systems look like we consider the change of variable $r = e^s$ and use $p(s)$ to denote $p(e^s)$ for forms of type I and we use $l(s), u(s)$ to denote $l(e^s)$ and $u(e^s)$ respectively for forms of type II. With this notation we have for example

$$p' = e^{-s} \partial_s p, \quad p'' = e^{-2s} \left( \partial^2_s p - \partial_s p \right).$$

On forms of type I, $\Box \xi$ is given by

$$\boxed{\Box \xi = e^{-2s} \left( \partial^2_s p + (n-4) \partial_s p - \mu p \right) \psi,}$$

while on forms of type II, $\Box \xi$ takes the form

$$\boxed{\Box \xi = e^{-2s} \left( 2 \partial^2_s l + 2(n-4) \partial_s l - 4(n-2 + \frac{\nu}{4}) l - \nu \partial_s u + 4\nu u \right) \phi ds + e^{-2s} \left( \partial^2_s u + (n-4) \partial_s u - 2\nu u + \partial_s l + nl \right) \tilde{d} \phi.}$$

The solutions of (4.18) and (4.19) are given by the following

$$\alpha = \frac{4-n}{2}, \quad \theta = \sqrt{\alpha^2 + \mu}, \quad a^\pm = \alpha \pm \theta.$$

All solutions of $\Box \xi = 0$ of type I are given by

$$\xi = r^{a^\pm} \psi.$$

In this case, since $r\psi$ is radially parallel, we see that the order of growth of $\xi$ is $a^\pm - 1$. For solutions $\xi$ of type II, we set

$$\beta = \frac{2-n}{2}, \quad \omega = \sqrt{\beta^2 + \nu}, \quad b^\pm = \beta \pm \omega,$$

and then $\xi$ is either of the form

$$r^{b^\pm} \tilde{d} \phi + b^\pm r^{b^\pm - 1} \phi dr,$$

or of the form

$$2r^{b^\pm + 2} \tilde{d} \phi + b^\pm r^{b^\pm + 1} \phi dr.$$

See [CT94, Section 2] for more details. The above computations motivate the following definition

**Definition 4.1.** The set $E$ of all numbers $a^\pm - 1, b^\pm - 1$ or $b^\pm + 1$ with $a^\pm, b^\pm$ defined by (4.20), (4.22), is called the set of *exceptional values for* $\Box$. If $\gamma \in \mathbb{R} \setminus E$ then $\gamma$ is said to be *nonexceptional*.

**Remark 4.2.** All elements in $E$ are integers, in fact, computing the eigenvalues of $\tilde{\Delta}_H$ on 1-forms in $S^{n-1}$ as in [Fol89, Theorem C], one can prove that all numbers in (4.21) have the form

$$a^\pm = -\frac{(n-4)}{2} \pm \frac{1}{2} (n-2+2j) \text{ for } j = 1, 2, \ldots,$$
and all numbers in (4.22) have the form
\[ b_j^\pm = -\frac{(n - 2)}{2} \pm \frac{1}{2} (n - 2 + 2j) \quad \text{for} \quad j = 0, 1, 2, \ldots \]

An important property of \( \Box \) is

**Proposition 4.3.** On \( \mathbb{R}^n \setminus \{0\} \), if \( \gamma \) is nonexceptional then \( \Box : W^{2,p,0,1}_\gamma \rightarrow W^{0,p,0,1}_{\gamma-2} \) is an isomorphism with bounded inverse.

**Proof.** Compare [Bar86, Theorem 1.7]. \( \Box \)

Finally, note that from (4.25) and (4.26) it follows that 1 is an exceptional value for \( \Box \), which means that there are elements in the kernel of \( \Box \) with linear growth, i.e., forms \( \xi \) with \( \Box \xi = 0 \) satisfying \( \xi = r\eta \) where \( \eta \) is radially parallel.

All 1-forms of type I in the kernel of \( \Box \) which have linear growth have the form
\[ \xi = r^{2\psi}, \]
where \( \psi \) is dual to a Killing field in \( S^{n-1} \). For forms of type II, all solutions of \( \Box \xi = 0 \) that have linear growth correspond to the eigenvalues \( \nu = 0, 2n \) of \( \tilde{\Delta}_H \) on functions, moreover, in that case the solution corresponding to \( \nu = 0 \) is
\[ \xi = r \phi \psi + r^{2\phi} \]
and the solution corresponding to \( \nu = 2n \) is
\[ \xi = 2r\phi \psi + r^{2\phi} \phi. \]

Note that the forms in (4.27) and (4.29) have linear growth because \( r\psi \) and \( r\phi \phi \) are radially parallel.

4.2. **A modified \( \Box \) operator.** Given \( t \neq 0 \), let \( \Box_t \) be the modified operator
\[ \Box_t \xi \equiv \delta_t L_\xi g_0 = \Box \xi - t i_{r^{-1}} \frac{\bar{\gamma}}{dr} L_\xi g_0. \]

In order to compute \( \Box_t \) for \( \xi \) as in (4.12) we start by noting that
\[ i_{r^{-1}} \frac{\bar{\gamma}}{dr} L_\xi g_0 = 2r^{-1} f' dr + r^{-1} \left( \omega' - 2r^{-1} \omega + \tilde{d}f \right). \]

The modified operator \( \Box_t \) is then computed to be
\[ \Box_t \xi = \left( 2f'' + 2(n - 1 - t)r^{-1} f' + r^{-2} \left( -\tilde{\Delta}_H f - 2(n - 1) f \right) \right) dr 
+ r^{-2} \left( \tilde{\delta} \omega' - 4r^{-3} \tilde{\delta} \omega \right) dr + (n - 3 - t)r^{-1} \omega' 
+ r^{-2} \left( -\tilde{\Delta}_H \omega + 2t \omega + \tilde{d} \tilde{\omega} \right) + \tilde{d} f' + r^{-1} (n + 1 - t) \tilde{d} f. \]

Using again the change of variable \( r = e^s \) and the notation in Section 4.1 we see that \( \Box_t \) for \( \xi \) a 1-form of type I is given by
\[ \Box_t \xi = e^{-2s} \left( \partial_s^2 p + (n - 4 - t) \partial_s p - (\mu - 2t)p \right) \psi. \]
On forms of type II, $\Box_t \xi$ is given by

\begin{equation}
(4.32) \quad \Box_t \xi = e^{-2s} \left( 2\partial^2_t l + 2(n - 4 - t)\partial_t l - 4(n - 2 - t)l - 2\nu^\prime l + 4nu \right) \phi ds + e^{-2s} \left( \partial^2_u u + (n - 4 - t)\partial_t u - 2(\nu - t)u + \partial_t l + (n - t)l \right) d\phi.
\end{equation}

We conclude that in these cases, the system

\begin{equation}
(4.33) \quad \Box_t \xi = 0,
\end{equation}

reduces again to a constant coefficient system of ordinary differential equations. As in the discussion at the end of Subsection 4.1, we are interested in solutions of \((4.33)\) which are essentially linear, i.e., solutions $\xi$ that satisfy for every $0 < \gamma < 1$,

\begin{equation}
(4.34) \quad \xi = O(r^{1+\gamma}) \text{ as } r \to \infty, \quad \xi = O(r^{1-\gamma}) \text{ as } r \to 0.
\end{equation}

For these solutions we have

**Proposition 4.4.** There exists $t_0 > 0$ such that for $0 < |t| < t_0$ there is a number $\gamma_0(t)$ with $0 < \gamma_0(t) < 1$ such that if $\xi$ is a 1-form in the kernel of $\Box_t$ satisfying \((4.34)\) for some $0 < \gamma < \gamma_0(t)$ then $\xi$ is dual to a Killing vector field in $\mathbb{R}^n$.

**Proof.** Since $\Box_t \xi = \delta_t L_\xi g_0$, it follows that a 1-form which is dual to a Killing field in $\mathbb{R}^n$ is in the kernel of $\Box_t$ regardless of the gauge condition. Next, since the general solution may be written as an infinite sum of 1-forms of Type I and Type II, it suffices to prove the proposition for 1-forms of either type. When $t = 0$ and as pointed out in Section 4.1, all solutions of \((4.33)\) in separated variables which satisfy \((4.34)\) correspond to the eigenvalues $\mu = 2(n - 2)$ on forms of type I and $\nu = 0, 2n$ on forms of type II. For $t \neq 0$ small, the growth rates are small perturbations of the rates $a^\pm_j - 1$ and $b^\pm_j - 1$ with $a_j^\pm$, $b_j^\pm$ given by \((4.25)\) and \((4.26)\), moreover the rates corresponding to eigenvalues $\mu > 2n$ or $\nu > 2n$ have real parts bounded away from 1 and hence, for these solutions the proposition follows. It only remains to consider the kernel corresponding to the eigenvalues $\mu = 2(n - 2)$ and $\nu = 0, 2n$. From [Fol89, Theorem C], all eigenvalues $\mu$ corresponding to 1-forms of type I are

\begin{equation}
(4.35) \quad \mu = (j + 1)(j + n - 3) \text{ for } j = 1, 2, \ldots,
\end{equation}

and then from \((4.31)\) all solutions of $\Box_t \xi = 0$ with $\xi$ a 1-form of type I can be written as

\begin{equation}
(4.36) \quad \xi = r c_j^\pm(t)^{-1} (r \psi_j),
\end{equation}

where

\begin{equation}
(4.37) \quad c_j^\pm(t) = \frac{-(n - 4 - t) \pm \sqrt{(n - 2 + 2j)^2 - 2nt + t^2}}{2},
\end{equation}

and $\psi_j$ is a 1-form with eigenvalue $(j+1)(j+n-3)$. It follows that if $t \neq 0$ is sufficiently small, all solutions given by \((4.37)\) are such that $Re \left( c_j^\pm(t) \right) - 1$ is bounded away from 1 except for $c_j^\pm(t) - 1$ which equals 1 for any $t$. In this case $r^2 \psi_1$ is dual to a Killing field in $\mathbb{R}^n$. For those 1-forms of type II corresponding to the eigenvalues $\nu = 0, 2n$,
the growth rates are strictly bounded away from 1 for \( t \neq 0 \) sufficiently small; the proof is similar and is omitted. □

5. Existence of divergence-free gauges

In order to construct divergence-free gauges, we use the approach in [CT94] which consists in using the ALE of order 0 condition to initially prove that we can fix a gauge such that the modified divergence-free condition \( \delta_t h = 0 \) for \( t \neq 0 \) is satisfied. One is then interested in the \( \delta_t \)-free kernel of the modified operator \( \mathcal{P}_t^{(k)} \). This kernel can be classified into three types:

1. Growth solutions, i.e., solutions that are \( O(r^\beta') \) on \( \mathbb{R}^n\setminus\{0\} \) for some \( \beta' > 0 \),
2. Decay solutions, i.e., solutions that are \( O(r^{-\beta'}) \) on \( \mathbb{R}^n\setminus\{0\} \) for some \( \beta' > 0 \),
3. “Degenerate” solutions, i.e., solutions that are \( O(r^\gamma) \) as \( r \to \infty \) and \( O(r^{-\gamma}) \) as \( r \to 0 \) for some \( \gamma > 0 \) sufficiently small (depending upon \( t \)).

The main step is to prove that solutions with the behavior described in (3) do not occur (see Proposition 5.4). We can then prove a growth estimate for solutions of the linear elliptic equation \( \mathcal{P}_t^{(k)} h = 0 \) (see Lemma 5.3) which we call the Three Annulus Lemma. The next step is to use scaling properties of the nonlinear system (3.28), and elliptic estimates to prove a nonlinear version of the Three Annulus Lemma (Lemma 5.8), so that the behavior of solutions of (3.28) can be modeled after the behavior of solutions of the linearized equation. Consequently, we can use the nonlinear Three Annulus Lemma and the ALE of order zero condition to rule out the behavior described in (1) and (3) for solutions of the nonlinear equation. It follows that the only possible behavior at infinity for solutions of the nonlinear equation in the \( \delta_t \)-gauge is that of decay solutions in (2), which yields a gauge where the metric \( g \) is ALE of positive order (see Corollary 5.12). With this improvement, one can easily construct a global divergence-free gauge, see Proposition 5.13.

The Three Annulus Lemma was introduced in [Sim85] and used in [CT94] for the Ricci-flat case. Even though our statement of the Three Annulus Lemma is very similar to that of [CT94], we base our proof on a result called Turan’s Lemma that we discuss in Appendix A. Our case is complicated by the fact that higher powers of log may enter into the asymptotic expansions since the system is of higher order.

5.1. The linearized equation in separated variables. We consider solutions on \( \mathbb{R}^n\setminus\{0\} \) of the system

\[
\mathcal{P}_t^{(k)} h = 0,
\]

for \( 1 \leq k \leq \frac{n}{2} - 1 \) if \( n \geq 4 \), or for \( k = 1 \) if \( n = 3 \), with \( \mathcal{P}_t^{(k)} \) given by (3.25). With notation as in [CT94], we write the general solution of \( \mathcal{P}_t^{(k)} h = 0 \) as an infinite sum of the form

\[
h = \sum_{j=0}^{\infty} \left( l_j \phi_j dr \otimes dr + k_j \tau_j \otimes dr + f_j B_j + p_j \phi_j r^2 g \right),
\]

where
(1) The functions $l_j, k_j, f_j$ and $p_j$ are radial, and the $(0,2)$ tensors $\phi_j dr \otimes dr$, $\tau_j \otimes dr$, $B_j$, and $\phi_j r^2 \tilde{g}$ are radially parallel.

(2) The components $\phi_j \in \Lambda^0(T^*S^{n-1})$, $\tau_j \in \Lambda^1(T^*S^{n-1})$, $B_j \in S^2_0(T^*S^{n-1})$ are eigenfunctions of the rough laplacian. Here $S^2_0(T^*S^{n-1})$ is the traceless component of $S^2(T^*S^{n-1})$.

(3) The set
\[
\{l_j \phi_j dr \otimes dr + k_j \tau_j \otimes dr + f_j B_j + p_j \phi_j r^2 \tilde{g}\}_{j},
\]
is orthogonal in the inner product $\langle \cdot, \cdot \rangle$ in $\mathbb{R}^8$.

It is clear that $\mathcal{P}_t^{(k)}$ preserves the expansion (5.2) in the sense that
\[
\mathcal{P}_t^{(k)} \left( l_j \phi_j dr \otimes dr + k_j \tau_j \otimes dr + f_j B_j + p_j \phi_j r^2 \tilde{g} \right) = F_j^{(1)} \phi_j dr \otimes dr + F_j^{(2)} \tau_j \otimes dr + F_j^{(3)} B_j + F_j^{(4)} \phi_j r^2 \tilde{g},
\]
where $F_j^{(c)} = F_j^{(c)}(l_j, k_j, f_j, p_j)$ for $c = 1, 2, 3, 4$. It follows that if $\mathcal{P}_t^{(k)} h = 0$ then for each $j$ we must have
\[
\mathcal{P}_t^{(k)} \left( l_j \phi_j dr \otimes dr + k_j \tau_j \otimes dr + f_j B_j + p_j \phi_j r^2 \tilde{g} \right) = 0,
\]
and
\[
F_j^{(c)}(l_j, k_j, f_j, p_j) \equiv 0, \text{ for } c = 1, 2, 3, 4.
\]
The system (5.4) is a linear system of ordinary differential equations which is homogeneous of order $2(k+1)$. Using the change of variable $r = e^{s}$, (5.4) can be written as a constant coefficient linear system of ordinary differential equations, in particular, for every $j = 0, 1, \ldots$, the system (5.4) reduces to a first order constant coefficient linear system of the form
\[
\dot{X} = M_j X,
\]
where $M_j$ is a matrix of order $8(k+1) \times 8(k+1)$. Let $\Phi_j$ is the characteristic polynomial of the matrix $A_j$ and suppose that we factor $\Phi_j$ as
\[
\Phi_j(z) = \prod_{a=1}^{m_j} (z - \zeta_{j,a})^{n_{j,a}},
\]
with
\[
\sum_{a=1}^{m_j} n_{j,a} = 8(k+1).
\]
Note that in general, the roots $\zeta_{j,a}$ depend on $t$ and $k$, however, for simplicity we omit this dependence in the notation above. Each of the functions $l_j, k_j, f_j, p_j$ may be expressed as a linear combination of functions of the form
\[
(\log(r))^{b} r^{k_{j,a}} \text{ with } b = 0, \ldots, n_{j,a} - 1.
\]
Fix $j$ and let $T_j^{(c)}$ with $c = 1, 2, 3, 4$, be the $(0,2)$ tensors
\begin{equation}
T_j^{(1)} = \phi_j dr \otimes dr, \quad T_j^{(2)} = \tau_j \boxtimes dr, \quad T_j^{(3)} = B_j, \quad T_j^{(4)} = \phi_j r^2 \tilde{g},
\end{equation}
so that we can write
\begin{equation}
l_j \phi_j dr \otimes dr + k_j \tau_j \boxtimes dr + f_j B_j + p_j \phi_j r^2 \tilde{g} = \sum_{c=1}^{4} q_{j,c} T_j^{(c)},
\end{equation}
where
\begin{equation}
q_{j,1} = l_j, \quad q_{j,2} = k_j, \quad q_{j,3} = f_j \quad \text{and} \quad q_{j,4} = p_j.
\end{equation}

It follows that if $h$ satisfies $P_k h = 0$ on $\mathbb{R}^n \setminus \{0\}$ then $h$ can be expanded as an infinite sum of the form
\begin{equation}
h = \sum_{j=0}^{\infty} \sum_{c=1}^{4} q_{j,c} T_j^{(c)},
\end{equation}
where
\begin{enumerate}
\item The $(0,2)$-tensors $T_j^{(c)}$ are radially parallel.
\item The set \{\(T_j^{(c)}\)\}_{j,c} is orthogonal with respect to the norm \langle \langle \cdot, \cdot \rangle \rangle.
\item The radial functions $q_{j,c}$ are linear combinations of functions of the form (5.8).
\end{enumerate}

Let us write the radial function $q_{j,c}(r)$ as
\begin{equation}
q_{j,c}(r) = \sum_{a=1}^{m_j} n_{j,a} \sum_{b=0}^{n_{j,a}-1} d_{a,b,c}(\log(r))^b r^{\zeta_{j,a}},
\end{equation}
where $d_{a,b,c}$ are complex numbers. From (5.13) we introduce the following sets
\begin{align}
A_j^+ &= \{1 \leq a \leq m_j : Re(\zeta_{j,a}) > 0\}; \\
A_j^- &= \{1 \leq a \leq m_j : Re(\zeta_{j,a}) < 0\}; \\
A_j^0 &= \{1 \leq a \leq m_j : Re(\zeta_{j,a}) = 0\}.
\end{align}

We will use $q_{j,c}^{\pm}$ to denote
\begin{equation}
q_{j,c}^{\pm}(r) = \sum_{a \in A_j^\pm} \sum_{b=0}^{n_{j,a}-1} d_{a,b,c}(\log(r))^b r^{\zeta_{j,a}},
\end{equation}
and $q_{j,c}^0$ will be used to denote
\begin{equation}
q_{j,c}^0(r) = \sum_{a \in A_j^0} \sum_{b=0}^{n_{j,a}-1} d_{a,b,c}(\log(r))^b r^{\zeta_{j,a}}.
\end{equation}

With this we have a decomposition of the form
\begin{equation}
h = h^+ + h^- + h^0,
\end{equation}
where

\begin{equation}
(5.20) \quad h^\pm = \sum_{j=0}^{\infty} \sum_{c=1}^{4} q_{j,c}^h T_j^{(c)}, \quad \text{and} \quad h^0 = \sum_{j=0}^{\infty} \sum_{c=1}^{4} q_{j,c}^0 T_j^{(c)}. \end{equation}

**Definition 5.1.** A solution \(h\) of (5.8) is a **degenerate solution of (5.1)** if \(h = h^0\).

It will also be important for us to consider for any nonnegative integer \(j\) the number

\begin{align*}
\beta_j &= \min \{|\text{Re}(\zeta_{j,a})| : a \in A_j^\pm\}.
\end{align*}

**5.2. Estimates for the linearized equation.** Consider a solution of (5.1) on an annulus, i.e., a solution of the problem

\begin{equation}
(5.22) \quad P_t^{(k)} h = 0 \text{ on } A_{a,b}(0),
\end{equation}

where \(0 < a < b\). Note that since \(A_{a,b}(0)\) and \(\mathbb{R}^n \setminus \{0\}\) have the same cross section, if \(h\) is a solution of (5.22) then we can repeat the analysis in Subsection 5.1 to decompose \(h\) as

\begin{equation}
(5.26) \quad h = h^+ + h^- + h_0 \text{ on } A_{a,L^2a}(0).
\end{equation}

Then there exists \(L_0 = L_0(\beta, \beta') > 1\) such that if \(L > L_0\), then

\begin{equation}
(5.27) \quad |||h^+|||_{L^2a,L^2a} \geq L^{\beta'} |||h^+|||_{a,L^2a},
\end{equation}

and

\begin{equation}
(5.28) \quad |||h^-|||_{L^2a,L^2a} \leq L^{-\beta'} |||h^-|||_{a,L^2a}.
\end{equation}

**Proof.** By the scale invariance of the norms \(||| \cdot |||\), it suffices to prove the lemma for \(a = 1\). The proof is completed in Appendix A using Turan’s Lemma. \(\square\)
We note that we are only able to prove (5.27) and (5.28) for $0 < \beta' < \frac{1}{2} \beta$ and not for $0 < \beta' < \beta$ as in [CT94]. However, the estimates (5.27) and (5.28) are sufficient for our purpose. Next, we use this to prove

**Lemma 5.3** (Three Annulus Lemma). Let $a > 0$, $L > 1$ and suppose that $h$ is a solution of (5.1) in $A_{L,a}((0)$ for $t$ sufficiently small. Suppose in addition that in the decomposition (5.26), $h_0 \equiv 0$. For any $0 < \beta' < \frac{1}{2} \beta$, there exists $L_0 = L_0(\beta, \beta')$ with $L_0 > 1$ such that for $L > L_0$, if

\[(5.29) \quad |||h|||_{aL,aL^2} \geq L^{\beta'}|||h|||_{a,aL},\]

then

\[(5.30) \quad |||h|||_{aL^2,aL^3} \geq L^{\beta'}|||h|||_{aL,aL^2},\]

and if

\[(5.31) \quad |||h|||_{aL^2,aL^3} \leq L^{-\beta'}|||h|||_{aL,aL^2},\]

then

\[(5.32) \quad |||h|||_{aL,aL^2} \leq L^{-\beta'}|||h|||_{a,aL}.\]

Moreover, at least one of (5.30), (5.32) holds (whether or not at least one of (5.29), (5.32) holds).

**Proof.** By scaling properties of the norms $||| \cdot |||$, it suffices to prove the lemma for the case $a = 1$. Suppose that (5.29) holds. From the decomposition $h = h^+ + h^-$ and the Cauchy-Schwarz inequality we clearly have

\[(5.33) \quad |||h|||_{L,L^2}^2 \leq 2 (|||h^+|||_{L,L^2}^2 + |||h^-|||_{L,L^2}^2),\]

and then

\[(5.34) \quad 2 (|||h^+|||_{L,L^2}^2 + |||h^-|||_{L,L^2}^2) \geq L^{2\beta'} (|||h^+|||_{1,L}^2 + |||h^-|||_{1,L}^2 + 2 \langle\langle h^+, h^- \rangle\rangle_{1,L}).\]

Here $\langle\langle \cdot, \cdot \rangle\rangle$ is the inner product associated to the norm $||| \cdot |||$. From Lemma 5.2 for $L$ large enough (depending on $\beta$ and $\beta'$) we can estimate $\langle\langle h^+, h^- \rangle\rangle_{1,L}$ in (5.34) as

\[(5.35) \quad \langle\langle h^+, h^- \rangle\rangle_{1,L} \geq -|||h^+|||_{1,L}|||h^-|||_{1,L} \geq -L^{-\beta'}|||h^+|||_{L,L^2}|||h^-|||_{1,L},\]

and then for a fixed $0 < \epsilon < \frac{1}{2}$, there exists a positive constant $c(\epsilon)$ such that (5.34) and (5.35) imply

\[(5.36) \quad L^{2\beta'} (|||h^+|||_{1,L}^2 + |||h^-|||_{1,L}^2 - 2\epsilon|||h^-|||_{1,L}^2 - 2c(\epsilon)L^{-2\beta'}|||h^+|||_{L,L^2}^2 - 2\epsilon|||h^-|||_{L,L^2}^2) \leq 2 (|||h^+|||_{L,L^2}^2 + |||h^-|||_{L,L^2}^2),\]

and hence

\[(5.37) \quad 2(1 + c(\epsilon)) (|||h^+|||_{L,L^2}^2 + |||h^-|||_{L,L^2}^2) \geq (1 - 2\epsilon)L^{2\beta'} (|||h^+|||_{1,L}^2 + |||h^-|||_{1,L}^2).\]
We have shown that for fixed $0 < \epsilon < \frac{1}{2}$ there exists a positive constant $q(\epsilon)$ such that

\begin{equation}
(|||h^+|||^2_{L^2, L^2} + |||h^-|||^2_{L^2, L^2}) \geq q(\epsilon)L^{2\beta'} (|||h^+|||^2_{L^2, L^2} + |||h^-|||^2_{L^2, L^2}).
\end{equation}

Combining Lemma 5.2 with (5.38) we obtain

\begin{equation}
(\beta' < \beta'' < \frac{1}{2}\beta, \text{ there exists } L_0 = L_0(\beta, \beta'') > 0 \text{ such that if } L > L_0 \text{ then}
\end{equation}

\begin{equation}
(|||h|||^2_{L^2, L^3}) \geq (1 - 2\epsilon)L^{2\beta''} |||h^+|||^2_{L^2, L^2} - 2c(\epsilon)|||h^-|||^2_{L^2, L^2},
\end{equation}

If $L$ is large enough so that $L^{-2\beta''} < \frac{1}{2}$ we have from (5.40)

\begin{equation}
L^{2\beta''} |||h^+|||^2_{L^2, L^2} \geq L^{2\beta''} \left(\frac{1}{2} + L^{-2\beta''}\right) |||h^+|||^2_{L^2, L^2} \geq \frac{1}{2}L^{2\beta''} |||h^+|||^2_{L^2, L^2} + |||h^-|||^2_{L^2, L^2}
\end{equation}

\begin{equation}
\geq \frac{1}{2}L^{2\beta''} |||h^+|||^2_{L^2, L^2} + \left(q(\epsilon)L^{4\beta'} - 1\right) |||h^-|||^2_{L^2, L^2}.
\end{equation}

Finally, from (5.42) and (5.43) it follows that

\begin{equation}
|||h|||^2_{L^2, L^3} \geq v(\epsilon, L, \beta', \beta'')L^{2\beta'} |||h^+|||^2_{L^2, L^2} + w(\epsilon, L, \beta', \beta'')L^{2\beta'} |||h^-|||^2_{L^2, L^2},
\end{equation}

where

\begin{equation}
v(\epsilon, L, \beta', \beta'') = \frac{1}{2}(1 - 2\epsilon)L^{2(\beta'' - \beta')},
\end{equation}

and

\begin{equation}
w(\epsilon, L, \beta', \beta'') = (1 - 2\epsilon)(q(\epsilon)L^{2\beta'} - L^{-2\beta'}) - 2c(\epsilon)L^{-2(\beta' + \beta')}.
\end{equation}

It is clear that we can choose $L$ large enough so that $v(\epsilon, L, \beta', \beta''), w(\epsilon, L, \beta', \beta'') \geq 2$ and then

\begin{equation}
|||h|||^2_{L^2, L^3} \geq 2L^{2\beta'} (|||h^+|||^2_{L^2, L^2} + |||h^-|||^2_{L^2, L^2}) \geq L^{2\beta'} |||h|||^2_{L^2, L^2},
\end{equation}

as needed. The proof for the case (5.31) is analogous. For the rest of the proposition, note that by the Cauchy-Schwarz inequality we must have either $|||h^+|||^2_{L^2, L^2} \geq \frac{1}{2}|||h|||^2_{L^2, L^2}$ or $|||h^-|||^2_{L^2, L^2} \geq \frac{1}{2}|||h|||^2_{L^2, L^2}$. If $|||h^+|||^2_{L^2, L^2} \geq \frac{1}{2}|||h|||^2_{L^2, L^2}$, then for fixed $0 < \epsilon < 1$ there exists $c_0(\epsilon) > 1$ such that

\begin{equation}
(1 - c_0(\epsilon))|||h^+|||^2_{L^2, L^2} + (1 - \epsilon)|||h^-|||^2_{L^2, L^2} \leq |||h|||^2_{L^2, L^2},
\end{equation}

and since $|||h|||^2_{L^2, L^2} \leq 4|||h^+|||^2_{L^2, L^2}$ we conclude that for some positive constant $c_1(\epsilon)$ we have

\begin{equation}
|||h^+|||^2_{L^2, L^2} \geq c_1(\epsilon)|||h^-|||^2_{L^2, L^2}.
\end{equation}
On the other hand, if we fix $0 < \epsilon < 1$ there exists a constant $c_2(\epsilon) > 0$ such that
\begin{equation}
\|h\|_{L^2,L^3}^2 \geq (1-\epsilon)\|h^+\|_{L^2,L^3}^2 - c_2(\epsilon)\|h^-\|_{L^2,L^3}^2,
\end{equation}
and from Lemma 5.2 for any $\beta' < \beta'' < \frac{1}{2}\beta$ there exists $L_0 = L_0(\beta,\beta'')$ such that if
$L > L_0$ then
\begin{equation}
\|h\|_{L^2,L^3}^2 \geq (1-\epsilon)L^{2\beta''}\|h^+\|_{L^2,L^2}^2 - c_2(\epsilon)L^{-2\beta''}\|h^-\|_{L^2,L^2}^2,
\end{equation}
and from (5.49) we have
\begin{equation}
\|h\|_{L^2,L^3}^2 \geq C(\epsilon, L, \beta', \beta'')L^{2\beta'}\|h^+\|_{L^2,L^2}^2,
\end{equation}
where
\begin{equation}
C(\epsilon, L, \beta', \beta'') = (1-\epsilon)L^{2(\beta''-\beta')} - \frac{c_2(\epsilon)}{c_1(\epsilon)}L^{-2(\beta'+\beta')}.
\end{equation}
If we choose $L$ large enough so that $C(\epsilon, L, \beta', \beta'') \geq 4$ we obtain
\begin{equation}
\|h\|_{L^2,L^3}^2 \geq 4L^{2\beta'}\|h^+\|_{L^2,L^2}^2 \geq L^{2\beta'}\|h\|_{L^2,L^2}^2,
\end{equation}
as needed. In a similar way we can show that if $\|h^-\|_{L^2,L^2} \geq \frac{1}{2}\|h\|_{1,L}^2$, then we have the inequality $\|h\|_{L^2,L^2}^2 \leq L^{-2\beta'}\|h\|_{1,L}^2$, which completes the proof. \hfill \Box

5.3. Degenerate solutions of the linearized equations. We now turn our attention to degenerate solutions of (5.1). If $t = 0$, then constants are non-trivial degenerate solutions, which are also divergence-free. However, these are not $\delta_t$-free for $t \neq 0$. The main result of this section is that there are in fact no degenerate solutions of (5.1) for all $t$ nonzero and sufficiently small:

**Proposition 5.4.** There exists $t_0 > 0$ such that if $0 < |t| < t_0$ there are no degenerate solutions of (5.1) subject to $\delta_t h = 0$ on any annulus $A_{c,d}(0)$. In particular, for $t \neq 0$ sufficiently small, Lemma 5.3 holds.

**Proof.** We only need to consider the case that $h$ is a finite sum in (5.12). In this case, $h$ extends to a solution of $P^{(k)}_t h = 0$ on $\mathbb{R}^n \setminus \{0\}$ subject to $\delta_t h = 0$. Let $\rho > 0$, let $t_0, t, \gamma_0(t)$ be as in Proposition 4.4 and let $\varphi$ be a $C^\infty$ function such that $\varphi|_{B_\rho(0)} \equiv 0$ and $\varphi|_{\mathbb{R}^n \setminus B_{2\rho}(0)} \equiv 1$. Choose $p > n$ and a number $0 < \gamma < \gamma_0(t) < 1$, then $\varphi \delta h \in W^{\gamma,p,0,1}_{-\gamma}$ and $(1-\varphi) \delta h \in W^{\gamma,p,0,1}_{-\gamma-1}$. Since $\gamma$ is nonexceptional it follows from Proposition 4.3 that there exists $X_1 \in W^{2,p,0,1}_{\gamma+1} \subset W^{2,p,0,1}_{\gamma+1}$ such that $\Box X_1 = \varphi \delta h$ and $X_2 \in W^{2,p,0,1}_{-\gamma+1}$ such that $\Box X_2 = (1-\varphi) \delta h$, therefore
\begin{equation}
\Box (X_1 + X_2) = \delta h,
\end{equation}
moreover, there exists $h_0$ such that $\delta h_0 \equiv 0$, $\varphi h_0 \in W^{1,p,0,2}_{\gamma}$, $(1-\varphi) h_0 \in W^{1,p,0,2}_{-\gamma}$, and
\begin{equation}
h = L(X_1 + X_2) g_0 + h_0 \text{ on } \mathbb{R}^n \setminus \{0\}.
\end{equation}
Note that from the weighted Sobolev inequality (see [Bar86, Theorem 1.2]), we must have
\begin{equation}
\tag{5.57}
h_0 = O(r^{-\gamma}) \text{ as } r \to 0, \quad \text{and } h_0 = O(r^\gamma) \text{ as } r \to \infty.
\end{equation}

Let \( X = X_1 + X_2 \), from \( P_i^{(k)} h = 0 \) and \( \delta_i h = 0, h \) satisfies
\begin{equation}
\tag{5.58}
\Delta^{k-1} \left( \frac{1}{2} \Delta^2 h + \frac{1}{2} \nabla^2 \delta \delta h + \Delta \delta^* \delta h \right) = 0
\end{equation}
on \( \mathbb{R}^n \backslash \{0\} \). From diffeomorphism invariance, any Lie derivative \( L_X g_0 \) is in the kernel of the linearized operator, and similarly, \( R'_0 (L_X g_0) = 0 \). It follows that \( L_X g_0 \) satisfies (5.58) and so does \( h_0 \) on \( \mathbb{R}^n \backslash \{0\} \). From \( \delta h_0 = 0 \), \( h_0 \) satisfies
\begin{equation}
\tag{5.59}
\Delta^{k+1} h_0 = 0,
\end{equation}
so we can expand \( h_0 \) in terms of homogeneous solutions of (5.59) on \( \mathbb{R}^n \backslash \{0\} \) and by (5.57) and the proof of Proposition 2.2 we conclude that \( h_0 = \log(r) \cdot C + C' \) where \( C, C' \) are matrices whose components are constant, but since \( \delta h_0 \equiv 0 \) it follows from Proposition 2.2 that \( C = 0 \) and \( h_0 \) is constant in \( \mathbb{R}^n \), in particular \( h_0 \) is a Lie derivative. We can now write \( h = L_Y g_0 \) where \( Y \) is a solution of \( \Box_i Y = \delta_i L_Y g_0 \equiv 0 \) and clearly \( Y \) is essentially linear in the sense of (4.34) for some \( \gamma \) with \( 0 < \gamma < \gamma_0(t) \). But from Proposition 4.4, we know that for \( t \neq 0 \) sufficiently small, if any such solution \( Y \) is non-zero then it must be dual to a Killing field which shows that \( h \equiv 0 \) as needed. In the case \( n = 3 \), all solutions \( \Delta^2 h_0 = 0 \) satisfying (5.57) are of the form \( h_0 = C + h_1 \) where the components of \( C \) are constant and the components of \( h_1 \) are spherical harmonics of order 1, however, if \( \delta h_0 \equiv 0 \) then \( h_1 \equiv 0 \) as seen in the proof of Proposition 2.2. \( \square \)

**Remark 5.5.** The argument in [Str10, Corollary 3.7] is incomplete since only radially parallel solutions are ruled out there. One must moreover rule out degenerate solutions (those with oscillatory behavior and possibly times a power of log) which are not radially parallel.

### 5.4. Scaling and the nonlinear equation

In this subsection we prove the nonlinear version of the Three Annulus Lemma. We assume that \( (M^n, g) \) is ALE of order 0, as in Definition 1.1. In the following we use the ALE coordinate system to transfer the problem to \( (\mathbb{R}^n \backslash B_\rho(0)) / \Gamma \). We have the following elliptic Schauder estimate for solutions of (3.28):

**Lemma 5.6.** Let \( 0 < a < d \) and let \( h \in \mathcal{T}^{2,0}_{m,\alpha,\beta}(A_{a,d}) \) be a solution of (3.28). There exists \( \chi > 0 \) such that if \( \|h\|_{\mathcal{T}^{2,0}_{m,\alpha,\beta}(A_{a,d}(0))} < \chi \), then for every \( b, c \) with \( a < b < c < d \) one has
\begin{equation}
\tag{5.60}
\|h\|_{\mathcal{T}^{2,0}_{m,\alpha,\beta}(A_{b,c}(0))} \leq C \|h\|_{A_{a,d}},
\end{equation}
with \( C = C(\lambda_1, \lambda_2, n, m, \alpha, b - a, d - c, t) \), where \( 0 < \lambda_1 \leq \lambda_2 \) are ellipticity constants of (3.28).
Proof. The result follows from standard interior elliptic regularity estimates, see for example [Eidelman 1969, Chapter II]. Note the leading order term is a power of the Laplacian, but lower order coefficients are negative powers of $r$. However, the Schauder estimate depends only on an appropriate weighted norm of the coefficients, which in this case is bounded, as one can easily verify.

The following scaling lemma will be used to reduce the nonlinear problem to the linear case.

Lemma 5.7. Let $\{h_i\}$ be a sequence of solutions of (3.28) satisfying
\begin{equation}
\|h_i\|_{T_0,2} < \chi_i, \tag{5.61}
\end{equation}
where $\{\chi_i\}$ is a sequence of positive numbers such that $\chi_i \to 0$. Suppose in addition that for some positive constant $C$ we have
\begin{equation}
|||h_i|||_{a,La} + |||h_i|||_{L^2,La} \leq C|||h_i|||_{La,L^2a}. \tag{5.62}
\end{equation}
Let $q_i = |||h_i|||_{La,L^2a}^{-1} h_i$, then on any annulus $A_{c,d}(0)$ with $a < c < d < L^3 a$, there exists a subsequence $q_{i_j}$ that converges in $T_0^{0,2}(A_{c,d}(0))$ with $\alpha' < \alpha$ to $\tilde{q}$ satisfying (5.1).

Proof. The sequence $\{q_i\}$ satisfies
\begin{equation}
|||q_i|||_{a,La} + |||q_i|||_{L^2,La} \leq C, \tag{5.63}
\end{equation}
in particular, if $c, d$ as in the statement, we have from Lemma 5.6 the inequality
\begin{equation}
|||q_i|||_{T_0^{0,2}(A_{c,d}(0))} \leq C', \tag{5.64}
\end{equation}
for some positive constant $C'$. By the Arzela-Ascoli theorem there exists a subsequence $\{q_{i_j}\}$ that converges in $T_0^{0,2}(A_{c,d}(0))$ with $\alpha' < \alpha$ to $\tilde{q}$. It only remains to prove that $\tilde{q}$ solves (5.1). For that purpose we write (3.28) as
\begin{equation}
P_t^{(k)}(c_i q_i) + R^{(k)}(c_i q_i, g_0) = 0, \tag{5.65}
\end{equation}
where $c_i = |||h_i|||_{La,L^2a}$. Note that $c_i \to 0$ as $i \to \infty$. From the estimate (5.64) and equation (3.27) it follows that
\begin{equation}
R^{(k)}(c_i q_i, g_0) = O(c_i^2) \text{ as } c_i \to 0, \tag{5.66}
\end{equation}
where the bound in the right-hand side of (5.66) is with respect to the norm in $T_0^{0,2}(A_{c,d}(0))$, so (5.65) takes the form
\begin{equation}
c_i P_t^{(k)}(q_i) = O(c_i^2) \text{ as } c_i \to 0. \tag{5.67}
\end{equation}
Passing to the subsequence $\{q_{i_j}\}$ we conclude that the limit $\tilde{q}$ satisfies $P_t^{(k)} \tilde{q} = 0$. □

Using Lemma 5.3 and Proposition 5.4 we have the following nonlinear version of the Three Annulus Lemma:
Lemma 5.8. Let \( \rho, t > 0 \) and let \( h \) be a solution of (3.28) on \( A_{\rho, \infty}(0) \) with \( \delta_t h = 0 \). Let \( \beta' > 0 \) and \( \delta_t > 0 \) as in Lemma 5.3 and let \( L, a > 0 \) be such that \( \delta_t \rho > \rho \) and \( L > L_0 \). There exist \( \chi = \chi(n, \lambda, \Lambda) > 0 \) so that if \( |h|_{T_{m, \alpha; 0}(A_{\rho, \infty}(0))} < \chi \), then if

\[
|||h|||_{L^2, L^2 a} \geq L^{\beta'} |||h|||_{a, L^2 a},
\]

then

\[
|||h|||_{L^2, L^2 a} \geq L^{\beta'} |||h|||_{L^2, L^2 a},
\]

and if

\[
|||h|||_{L^2, L^2 a} \leq L^{-\beta'} |||h|||_{L^2, L^2 a},
\]

then

\[
|||h|||_{L^2, L^2 a} \leq L^{-\beta'} |||h|||_{a, L^2 a}.
\]

Moreover, there exists \( t_0 > 0 \) such that if \( 0 < |t| < t_0 \) then at least one of (5.69), (5.71) must hold.

Proof. If any of the implications in the statement of the lemma fails we can use the rescaling construction in Lemma 5.7 to produce a solution of the linearized equation that contradicts Lemma 5.3. □

5.5. Global divergence-free gauges. From (1.1) and (1.2) it follows that \( \Psi^* g - g_0 \) satisfies the estimate

\[
| (\Psi^* g - g_0)_{(r, x)} |_{m, \alpha; 0} = o(1) \quad \text{as} \quad r \to \infty.
\]

for all \( m \geq 0 \).

Remark 5.9. We do not need to assume decay on derivatives of arbitrary order in (1.2). Our proof only requires that the estimate (1.2) be satisfied only for all multi-indices \( l \) such that \( 1 \leq |l| \leq 2(k + 1) + 1 \). Then (5.72) holds for all \( m \leq 2(k + 1) \).

With the decay in (5.72) we can prove the existence of \( \delta_t \)-free gauges on certain annuli by means of the implicit function theorem. The following result is contained in [CT94, Theorem 3.1] and does not depend on the metric \( g \) being \( \Omega^{(k)} \)-flat:

Proposition 5.10. Let \( t_0 \) be as in Proposition 4.4 and let \( t \) be such that \( 0 < |t| < t_0 \). There exists \( \chi = \chi(t, m) \) such that if \( (C(N^{n-1}), g_0) \) is a Ricci-flat cone and \( \tilde{g} \) is a metric on \( A_{c, \infty}(0) \subset C(N^{n-1}) \) such that

\[
\| \tilde{g} - g_0 \|_{T_{m, \alpha; 0}(A_{c, \infty}(0))} < \chi,
\]

then there exists a diffeomorphism \( \phi = \phi(\tilde{g}), \phi : A_{c, \infty}(0) \to A_{c, \infty}(0) \) such that

\[
\phi^* \tilde{g} \in T_{m, \alpha; 0}(A_{c, \infty}(0)),
\]

and

\[
\delta_t (\phi^* \tilde{g} - g_0) = 0.
\]
Moreover, if
\[(5.76)\quad \|\delta_t(g)\|_{T_{m,0}^{0,1}(A_{c,\infty}(p))} < \epsilon,\]
then
\[(5.77)\quad \|\phi^*g - \tilde{g}\|_{T_{m,0}^{0,2}(A_{c,\infty}(p))} < \delta(\epsilon),\]
where \(\delta(\epsilon) \searrow 0\) as \(\epsilon \to 0\).

In our application of this Lemma, we will simply let \(\tilde{g} = \Psi* g\), and \(g_0\) the flat metric on the Euclidean cone (recall Definition 5.1). We will then write \(h = \phi^* \tilde{g} - g_0\). From Proposition 5.10 and using that \(g\) is ALE of order 0 we have

**Lemma 5.11.** Let \(h\) and \(A_{c,\infty}(0)\) be as above, and let \(L_0 > 0\) be as in Lemma 5.3. For any \(L > L_0\) and any \(a > c\) we have \(\lim_{i \to \infty} |||h|||_{L^a,L^{i+1}a} = 0\), moreover, we have the inequality
\[(5.78)\quad |||h|||_{L^a,L^{i+1}a} \leq L^{-\beta'(i-i_0)} |||h|||_{L^a,L^{i_0+1}a} \quad \text{for all } i \geq i_0.\]

**Proof.** By Lemma 5.8 for any \(i > i_0\) we have either
\[(5.79)\quad |||h|||_{L^{i+1},L^{i+2}} \geq L^{\beta'} |||h|||_{L^i,L^{i+1}},\]
or
\[(5.80)\quad |||h|||_{L^i,L^{i+1}} \leq L^{-\beta'} |||h|||_{L^{i-1},L^i},\]
Suppose that we have (5.79), then we conclude again from Lemma 5.8 that we must have for any integer \(s \geq 1\) the inequality
\[(5.81)\quad |||h|||_{L^{i+s},L^{i+s+1}} \geq L^{\beta' s} |||h|||_{L^i,L^{i+1}}.\]
Since \(g\) is ALE of order 0, we can use (5.76) in Proposition 5.10 to conclude that for \(\epsilon > 0\) sufficiently small and \(s > 1\) large we must have
\[(5.82)\quad \|\Psi* g - g_0\|_{T_{m,0}^{0,2}(A_{L^{i+s},L^{i+s+1}(0)})} < \epsilon,\]
\[(5.83)\quad \|\phi^*\Psi* g - \Psi* g\|_{T_{m,0}^{0,2}(A_{L^{i+s},L^{i+s+1}(0)})} < \frac{1}{2},\]
and therefore
\[(5.84)\quad |||h|||_{T_{m,0}^{0,2}(A_{L^{i+s},L^{i+s+1}(0)})} < 1.\]
In particular, we must have
\[(5.85)\quad |||h|||_{L^{i+s},L^{i+s+1}} \leq c_n \log(L),\]
where \(c_n\) is a dimensional constant. It is clear that (5.81) contradicts (5.85) for \(s\) large. It follows that (5.80) holds and then, by Lemma 5.8 we must have (5.78). □

**Corollary 5.12.** If \(g\) is ALE of order 0, scalar flat and (extended) obstruction-flat or satisfies (1.9), then there exists an annulus of the form \(A_{c,\infty}(0)\) and a diffeomorphism \(\phi : A_{c,\infty}(0) \to A_{c,\infty}(0)\) such that \(\|\phi^*\Psi* g - g_0\|_{m,0} = O(r^{-\beta'})\) as \(r \to \infty\) and therefore, \(g\) is ALE of order \(\beta'\).
Proof. From Lemma 5.6, Proposition 5.10, and Lemma 5.11 there exists a constant $C > 0$ such that for any $(r, x) \in A_{c, \infty}(0)$ with $r$ sufficiently large we have

$$|(\phi^* \Psi g - g_0)_{(r, x)}|_{m, \alpha; 0} \leq Cr^{-\beta'},$$

so the claim follows. □

From [CT94, Sections 2 and 3] we have

**Proposition 5.13.** Suppose $g$ is a metric defined on $\mathbb{R}^n \setminus B_\rho(0)$ and satisfies

$$(5.86) \quad g - g_0 = O(r^{-\tau}) \text{ as } r \to \infty,$$

then there exists a diffeomorphism $\phi : \mathbb{R}^n \setminus B_\rho(0) \to \mathbb{R}^n \setminus B_\rho(0)$ such that $h = \phi^* g - g_0$ satisfies $h \in T_{m, \alpha; -\tau}(\mathbb{R}^n \setminus B_\rho(0))$ and $\delta_{g_0} h = 0$.

We summarize the results of this section in the following corollary

**Corollary 5.14.** If $g$ is ALE of order 0, scalar flat and (extended) obstruction-flat or satisfies (1.9), then there exists a diffeomorphism $\Phi : M \setminus K \to (\mathbb{R}^n \setminus B_\rho(0))/\Gamma$ for some $\rho > 0$ and a compact set $K \subset M$ such that $h = \Phi^* g - g_0$ satisfies $h \in T_{m, \alpha; -\tau}(\mathbb{R}^n \setminus B_\rho(0))$ and $\delta_{g_0} h = 0$.

6. Optimal ALE order

6.1. Weighted Sobolev spaces and $\Delta^{k+1}$. In this section we state some properties of the weighted Sobolev spaces introduced in Section 5.3 that will be useful to improve the decay estimate for the metric $g$ derived in Section 5.5. Throughout this section we will work only with $(0, 2)$ tensors so when we write $W_{m, p}^{\delta}$ we actually mean the space $W_{\delta}^{m, p, 0.2}$. We start by defining the set of exceptional values for $\Delta^{k+1}$.

**Definition 6.1.** A number $\delta \in \mathbb{R}$ is said to be exceptional for $\Delta^{k+1}$ if $\delta$ is in the set

$$(6.1) \quad E = \begin{cases} \{ j \in \mathbb{Z} : j \neq -1, -2, \ldots, 2(k + 1) - (n - 1) \} & \text{if } n > 2(k + 1) \\ \mathbb{Z} & \text{if } n = 2(k + 1). \end{cases}$$

We say that $\delta$ is nonexceptional if $\delta \in \mathbb{R} \setminus E$.

**Remark 6.2.** The exceptional values for $\Delta^{k+1}$ correspond to the growth rates of solutions of $\Delta^{k+1} h = 0$ on the complement of a ball, however, when $n = 2(k + 1)$, as observed in the proof of Proposition 2.2 there are solutions of $\Delta^{k+1} h = 0$ on $\mathbb{R}^n \setminus \{0\}$ that are $O(\log(r))$ as $r \to \infty$.

**Lemma 6.3.** If $\delta$ is nonexceptional, the map $\Delta^{k+1} : W_\delta^{2(k+1), p} \to W_{\delta-2(k+1)}^{0, p}$ is an isomorphism.

**Proof.** See [Bar86] Theorem 1.7. □
Lemma 6.4. Suppose that $h$ is defined on $\mathbb{R}^n \setminus B_\rho(0)$ and satisfies $h = O(r^\delta)$ as $r \to \infty$ and assume that $\Delta^{k+1}h = O(r^{\delta'-2(k+1)})$ as $r \to \infty$ with $\delta' < \delta$. Then, for any $\tau > 0$ such that $\delta' + \tau$ is nonexceptional there exists $h' \in W^{2(k+1),p}_{\delta'+\tau}$ and a ball $B_{\rho'}(0)$ such that
\begin{equation}
\Delta^{k+1}(h - h') = 0 \text{ on } \mathbb{R}^n \setminus B_{\rho'}(0),
\end{equation}
Furthermore, if $2(k+1) < n$, there exists an exceptional value $j \leq \max\{\delta, \delta' + \tau\}$ such that
\begin{equation}
h - h' = p_j + O(r^{j-1}) \text{ as } r \to \infty,
\end{equation}
where $p_j$ is homogeneous of degree $j$ and satisfies $\Delta^{(k+1)}(p_j) = 0$ on $\mathbb{R}^n \setminus \{0\}$. If $n = 2(k+1)$, we may also have
\begin{equation}
h - h' = A \cdot \log(r) + O(1) \text{ as } r \to \infty,
\end{equation}
where the components of $A$ are constant.

Proof. Let $\varphi$ be a function in $C^\infty(\mathbb{R}^n)$ such that $\varphi \equiv 0$ on $B_\rho(0)$ and $\varphi \equiv 1$ on $\mathbb{R}^n \setminus B_{2\rho}(0)$, then $\Delta^{k+1} (\varphi h) \in W^{0,p}_{\delta'-2(k+1)+\tau}$ for every $\tau > 0$. If we choose $\tau$ in such a way that $\delta' + \tau$ is nonexceptional, by Lemma 6.3 there exists $h' \in W^{2(k+1)}_{\delta'+\tau}$ such that
\begin{equation}
\Delta^{(k+1)}(h \varphi - h') = 0,
\end{equation}
and on $\mathbb{R}^n \setminus B_{2\rho}(0)$ we have (6.2). The expansion (6.3) follows from the expansion at infinity of solutions of $\Delta^{k+1}h = 0$ on $\mathbb{R}^n \setminus B_1(0)$ (compare with the proof of Proposition 5.4).

6.2. Optimal decay. Suppose that $(M^n, g)$ is ALE of order 0, scalar-flat, and either $\Omega^{(k)}$-flat or satisfies (1.9). Corollary 5.14 showed that $g$ is ALE of order $\beta'$ for some $\beta' > 0$. In the next proposition we obtain the optimal value for $\beta'$ as stated in Theorems 1.2 and 1.4.

Proposition 6.5. Let $h$ be as in Corollary 5.14. If $g$ is $\Omega^{(k)}$-flat or satisfies (1.9), then $h \in T_{m,\alpha;2k-n}^0$ and $g$ is ALE of order $n - 2k$.

Proof. Let us treat first the case $(k + 1) = \frac{n}{2}$. Since $\delta h = 0$, $h$ satisfies (3.28) with $t = 0$, i.e.
\begin{equation}
\frac{C_n}{2(n-2)} \Delta^{\frac{n}{2}} h = R^{(\frac{n}{2}-1)}(h, g_0),
\end{equation}
where $R^{(\frac{n}{2}-1)}(h, g_0)$ is given by equation (3.27). Since $h = O(r^{-\beta'})$ as $r \to \infty$, the least decaying terms in (3.27) are those terms containing $\nabla^{\alpha_1} h \nabla^{\alpha_2} h$ with $\alpha_1 + \alpha_2 = n$, so from (6.6) we obtain $\Delta^{\frac{n}{2}} h = O(r^{-2\beta'-n})$, and by Lemma (6.4), for any $\tau > 0$ such that $-2\beta' + \tau$ is nonexceptional there exists $h_1 \in W^{m,p}_{2\beta'+\tau}$ such that $\Delta^{\frac{n}{2}} h - h_1 = 0$ on the complement of some ball. From the weighted Sobolev inequality, if we take $p > n$ then $h_1 = O(r^{-2\beta'+\tau})$ as $r \to \infty$ and clearly we can assume that $-2\beta' + \tau < -\beta'$, so that both $h, h_1$ have pointwise decay at infinity but $h_1$ has a better decay at infinity.
than \( h \). By Lemma 6.4 and since \(-1\) is the least negative exceptional value for \( \Delta^2 \), the difference \( h - h_1 \) has an expansion at infinity of the form

\[
(6.7) \quad h - h_1 = F_1 + O(r^{-2}),
\]

where \( F_1 \) is a homogeneous solution of degree \(-1\) of \( \Delta^2 h = 0 \) on \( \mathbb{R}^n \setminus \{0\} \). From the proof of Proposition 2.2 any such \( F_1 \) has the form

\[
(6.8) \quad (F_1)_{ij}(x) = u_{ij} \left( \frac{x}{|x|^2} \right),
\]

where \( u_{ij} \) are linear functions. We now claim that on the complement of some ball, \( h \) satisfies

\[
(6.9) \quad h = F_1 + O(r^{-1-\epsilon}),
\]

for some \( \epsilon > 0 \). If \(-2\beta' < -1\) we can choose \( \tau > 0 \) sufficiently small so that \( h \) satisfies (6.9). If not, we have \( h = O(r^{-2\beta' + \tau}) \) with \(-2\beta' + \tau < -\beta'\) and again from (6.6) it follows that \( \Delta^2 h = O(r^{-4\beta' + 2\tau}) \) and we can argue as above to obtain \( h = O(r^{-\min(1,4\beta' - 2\tau)}) \) for some \( \tau' > 0 \) such that \(-4\beta' + \tau' < -2\beta' + \tau\). It is clear that we can use induction to obtain (6.9). Note that if \( \delta F_1 \) is not identically zero then \( \delta F_1 = O(r^{-2}) \) as \( r \to \infty \) but we do not have \( \delta F_1 = O(r^{-2-\epsilon}) \) as \( r \to \infty \) for any \( \epsilon > 0 \), however, by (6.9) one has

\[
(6.10) \quad \delta h = \delta F_1 + O(r^{-2-\epsilon}) \quad \text{as} \quad r \to \infty,
\]

therefore, from \( \delta h \equiv 0 \) it follows that \( \delta F_1 \equiv 0 \), and by Proposition 2.2 \( F_1 \equiv 0 \) which shows that \( h = O(r^{-\gamma}) \) with \( \gamma > 1 \). By (6.6) one obtains \( \Delta^2 h = O(r^{-2\gamma - n}) \) and repeating the argument used to obtain (6.9) we can show that on the complement of some ball, \( h \) has an expansion of the form

\[
(6.11) \quad h = F_2 + O(r^{-2-\epsilon}) \quad \text{as} \quad r \to \infty,
\]

where \( F_2 \) is homogeneous of degree \(-2\) and satisfies \( \Delta^2 F_2 \equiv 0 \) on \( \mathbb{R}^n \setminus \{0\} \) and \( \epsilon \) is some positive number. This proves that \( h = O(r^{-2}) \) as \( r \to \infty \) as needed. For the case \( 2(k + 1) < n \), the only difference with the previous proof is that we have to consider homogeneous solutions of \( \Delta^2 h \equiv 0 \) on \( \mathbb{R}^n \setminus \{0\} \) that decay like \( r^{2k+1-n} \) and like \( r^{2k+1-n} \), but as shown in the proof of Proposition 2.2 these solutions are not divergence-free unless they are identically zero.

7. Singularity removal theorems

In this section we present the proofs of Theorems 1.8 and 1.9.

Lemma 7.1. Let \( g = g_0 + h \) be a metric defined on \( B_\rho(0) \setminus \{0\} \) with constant scalar curvature, and assume that \( g \) is either \( \Omega^{(k)} \)-flat or satisfies (1.13). Suppose in addition that \( \delta h = 0 \) on \( B_\rho(0) \setminus \{0\} \). Then, on \( B_\rho(0) \setminus \{0\} \), \( h \) satisfies the equation

\[
(7.1) \quad \mathcal{P}_t^{(k)} h + \mathcal{R}^{(k)}(h, g_0) = 0,
\]

where \( \mathcal{P}_t^{(k)} \) and \( \mathcal{R}^{(k)}(h, g_0) \) have the same expressions as in (3.25) and (3.27) respectively. The operator \( \mathcal{P}_t^{(k)} \) is elliptic.
Proof. Suppose that \( R(g_0 + h) = c \) where \( c \) is a constant. If we also have \( \delta_h = 0 \) on \( B_\rho(0) \setminus \{0\} \), then from \( R(g_0 + h) - R(g_0) = c \) we conclude that \( h \) satisfies the equation
\[
\Delta \text{tr}(h) = -c + t \delta_i r^{-1} \frac{\partial}{\partial r} h + F'(h, g_0),
\]
on \( B_\rho(0) \setminus \{0\} \). We now write the equation \( \Omega^{(k)}(g_0 + h) - \Omega^{(k)}(g_0) = 0 \) as
\[
0 = (\Omega^{(k)})_{g_0}'(h) + F^{(k)}(h, g_0),
\]
with \( F^{(k)}(h, g_0) \) as in (3.2). Using (2.1), we see that if we insert (7.2) into (7.3) then \( h \) satisfies (7.1) on \( B_\rho(0) \setminus \{0\} \). The rest of the claim follows easily, and the same argument works for (1.13).

Recalling the \( C^0 \)-orbifold condition as defined in Definition 1.6, we now assume that there exists a coordinate system around the origin such that
\[
g_{ij} = \delta_{ij} + o(1),
\]
(7.4)
\[
\partial^l g_{ij} = o(r^{-|l|}),
\]
(7.5)
for any multi-index \( l \) with \( |l| \geq 1 \) as \( r \to 0 \).

Remark 7.2. As in the ALE case (compare Remark 5.9), we do not need an assumption on derivatives of arbitrary order. If \( l \) in (7.5) only satisfies \( |l| \leq 2(k+1)+1 \), then we have \( |(g - g_0)_{(r,x)}|_{m,\alpha;0} = o(1) \) as \( r \to 0 \) for any \( m \leq 2(k+1) \), which is sufficient for our proof.

Next, we state the existence of a divergence-free gauge.

Lemma 7.3. Suppose that \( g \) defined on \( B_\rho(0) \) has constant scalar curvature and is (extended) obstruction-flat or satisfies (1.13). Suppose also that the origin is a \( C^0 \)-orbifold point for \( g \). Then for some \( \rho' < \rho \) there exists a diffeomorphism \( \phi : B_{\rho'}(0) \setminus \{0\} \to B_{\rho'}(0) \setminus \{0\} \) such that \( \delta_{g_0} \phi_\ast g = 0 \) on \( B_{\rho'}(0) \setminus \{0\} \). Moreover, there exists \( \sigma > 0 \) such that \( |\phi_\ast g - g_0| = O(r^\sigma) \) as \( r \to 0 \) and \( \partial^l \phi_\ast g = O(r^{\sigma - |l|}) \) for any multi-index \( l \) with \( |l| \geq 1 \).

Proof. This follows from a straightforward modification of the proof of Corollary 5.14.

We will also need the following

Lemma 7.4. If the components of \( h \in S^2(T^*\mathbb{R}^n) \) are linear functions then \( h = L_X g_0 \) for some quadratic vector field.

Proof. Let \( S^2_1 \) be the subspace of \( S^2(T^*\mathbb{R}^n) \) consisting of all elements whose components are linear functions. If \( h \in S^2_1 \), we can write the components of \( h \) as
\[
h_{ij}(x) = \sum_{l=1}^{n} A_{ijl} x_l,
\] (7.6)
where \( A_{ijk} \) is symmetric in \( i, j \) and therefore \( \dim(S^2_1) = \frac{n^2(n-1)}{2} \). On the other hand, if we let \( \Gamma^1_2 \) be the space of all vector fields whose components are functions which are homogeneous of degree 2, then any \( X \in \Gamma^1_2 \) can be written as \( X = X_i \frac{\partial}{\partial x^i} \), where

\[
X_i(\xi) = \sum_{l,m} a_{ilm} \xi_l \xi_m,
\]

with \( a_{ilm} \) symmetric in \( l, m \), and therefore \( \dim(\Gamma^1_2) = \frac{n^2(n-1)}{2} = \dim(S^2_1) \). Since there are no quadratic Killing vector fields, the map \( \mathcal{L} : \Gamma^1_2 \rightarrow S^2_1 \) defined by \( \mathcal{L}(X) = L_X g_0 \) is an isomorphism.

**Lemma 7.5.** Let \( X \) a vector field that is homogeneous of degree 2 and let \( K^*_X \) be the diffeomorphism generated by taking the flow of \( X \) to time 1 (which exists for \( r \) sufficiently small). If \( g_0 \) is the Euclidean metric we have

\[
K^*_X g_0 - L_X g_0 - g_0 = O(r^2) \text{ as } r \to 0.
\]

**Proof.** Let \( \phi_t \) be the flow of \( X \), then we have for any \( t > 0 \)

\[
\phi_t^* g_0 = g_0 + tL_X g_0 + E(t),
\]

where \( E(t) \) is an error term that can be estimated as

\[
|E(t)| \leq \frac{t^2}{2} \sup_{s \in [0,t]} \left| \left( \frac{\partial^2}{\partial s^2} \phi_s \right) \right|,
\]

here \( | \cdot | \) is the usual pointwise norm on \( S^2(T^*\mathbb{R}^n) \). In particular

\[
K^*_X g_0 = g_0 + L_X g_0 + E(1).
\]

Since \( X \) is homogeneous of degree 2, we have for any \( p \in \mathbb{R}^n \)

\[
\frac{\partial}{\partial t} |\phi_t(p)|^2 = 2\left( \frac{\partial}{\partial t} \phi_t(p), \phi_t(p) \right) \leq C|\phi_t(p)|^3,
\]

for some constant \( C > 0 \) that only depends on \( X \). Letting \( r \) denote the distance of \( p \) to the origin, it follows from \( (7.12) \) that for \( 0 < r < \frac{1}{C} \) and \( 0 \leq t \leq 1 \) we have the inequality

\[
|\phi_t(p)| \leq \frac{2r}{2-Crt},
\]

and then \( |\phi_t(p)| = O(r) \) as \( r \to 0 \). A similar argument shows that for all first order partial derivatives one has

\[
|\partial_t \phi_t| = O(1) \text{ as } r \to 0.
\]

Since \( g_0 \) is the Euclidean metric, we have

\[
(\phi^*_t g_0)_{ij} = \sum_{k,l} \partial_k(\phi_t)_i \partial_l(\phi_t)_j,
\]

and using the chain rule we can write schematically

\[
\frac{\partial^2}{\partial t^2} \phi^*_t g_0 = (\partial^2 X)(\phi_t) \cdot X(\phi_t) \cdot \partial \phi_t \cdot \partial \phi_t + (\partial X)(\phi_t) \cdot (\partial X)(\phi_t) \cdot \partial \phi_t \cdot \partial \phi_t,
\]
and by (7.14), (7.15) we conclude that for \( r \) sufficiently small
\[
(7.16) \quad \left| \frac{\partial^2}{\partial t^2} \phi_t g_0 \right| \leq C' r^2,
\]
for \( C' \) depending only on \( X \). By (7.10), (7.11) and (7.16) the result follows. \( \square \)

**Lemma 7.6.** Let \( g \) a metric defined on \( B_\rho(0) \) with a \( C^0 \)-orbifold point at the origin. Suppose that \( g \) has constant scalar curvature and is (extended) obstruction-flat or satisfies (1.13) on \( B_\rho(0) \setminus \{0\} \). Then there exists a change of coordinates \( \tilde{\phi} \), defined in some small neighborhood around the origin, such that \( \tilde{\phi}_* g \) satisfies
\[
(7.17) \quad g = g_0 + O(|x|^2),
\]
\[
(7.18) \quad \partial_t^l g = O(|x|^{2-|l|}),
\]
for any multi-index \( l \) with \( |l| \geq 1 \) as \( |x| \to 0 \).

**Proof.** Let \( \phi \) be the gauge given by Lemma 7.3 and if we take \( h = \phi_* g - g_0 \) we obtain
\[
(7.19) \quad \frac{c_{n,k}}{2(n-2)} \Delta^{k+1} h = \mathcal{R}^{(k)}(h, g_0).
\]
Since \( h = O(r^\sigma) \) as \( r \to 0 \) we have \( \Delta^{k+1} h = O(r^{2\sigma - 2(k+1)}) \) as \( r \to 0 \) and as in the proof of Proposition 6.5 for \( p > n \) and for \( \tau > 0 \) such that \( 2\sigma - \tau \) is nonexceptional and positive, there exists \( h' \in W^{m,p}_{2\sigma - \tau} \) such that \( \Delta^{k+1} (h - h') = 0 \) on \( B_\rho(0) \setminus \{0\} \). Since both \( h, h' \) are \( o(1) \) as \( r \to 0 \) we conclude that
\[
(7.20) \quad h - h' = G_1 + O(r^2) \text{ as } r \to 0,
\]
and the components of \( G_1 \) are linear functions. As in Proposition 6.5 we can use induction to show that \( h \) satisfies
\[
(7.21) \quad h = G_1 + O(r^{1+\epsilon}) \text{ as } r \to 0,
\]
for some \( \epsilon > 0 \). The strategy for proving (7.17), (7.18) is slightly different to that used to prove Proposition 6.5 but is still based on an argument used in [C194]. From \( \delta h = 0 \) on \( B_\rho(0) \setminus \{0\} \), it follows that \( \delta G_1 = 0 \) and by Lemma 7.4 \( G_1 = L_X g_0 \) for some vector field \( X \) such that \( X(p) \) is homogeneous of degree 2 in \( p \). Assume that \( \rho' \) is sufficiently small so that \( K_X \), the diffeomorphism obtained by taking the flow of \( X \) to time 1, is defined. By Lemma 7.5
\[
(7.22) \quad K^*_X g_0 - g_0 - L_X g_0 = O(r^2) \text{ as } r \to 0,
\]
and from
\[
(7.23) \quad (\phi_* g - K^*_X g_0) + (K^*_X g_0 - g_0 - L_X g_0) = O(r^{\min\{2,1+\epsilon\}}) \text{ as } r \to 0,
\]
we conclude that
\[
(7.24) \quad K^*_X \phi_* g - g_0 = O(r^{\min\{1+\epsilon,2\}}) \text{ as } r \to 0.
\]
As in Corollary 5.14 we can find a diffeomorphism \( \phi' \) defined on a smaller ball such that
\[
(7.25) \quad h' = \phi'_* K^*_X \phi_* g - g_0.
\]
is divergence-free and \( h' = O(r_{\min(1+\epsilon,2)}) \) as \( r \to 0 \). With this new \( h' \) we argue again as in the proof of Proposition 6.5 to obtain (7.17) and (7.18) as needed. \( \square \)

**Lemma 7.7.** In the coordinate system constructed in Lemma 7.6 we have

\[
(7.26) \quad \nabla^l Rm = \partial^l Rm + O(1) \quad \text{as} \quad r \to 0,
\]

for any multi-index \( l \) with \(|l| \geq 1\) and

\[
(7.27) \quad \Delta^m g \Ric(g) = \Delta^m g_0 \Ric(g) + O(r^{-2(m-1)}) \quad \text{as} \quad r \to 0.
\]

Proof. To show (7.26), we consider first the case \(|l| = 1\) and we write \( \nabla Rm \) schematically as

\[
(7.28) \quad \nabla Rm = \partial Rm + \Gamma^* Rm,
\]

and note that the terms \( \Gamma^* Rm \) are \( O(1) \) as \( r \to 0 \). The general case follows easily by induction. For (7.27) we start also with the case \( m = 1 \) and write

\[
(7.29) \quad \Delta g \Ric = g^{-1} \Delta \nabla \Ric = g^{-1} \nabla (\partial \Ric + \Gamma \ast \Ric)
\]

\[
(7.30) \quad \Delta g_{0 \ast \Ric} \ast (g^{kl} - \delta^{kl}) \partial_{kl} \Ric_{ij}.
\]

The terms \((g^{kl} - \delta^{kl}) \partial_{kl} \Ric_{ij}, \partial \Gamma \ast \Ric, \text{ and } \Gamma \ast \partial \Ric \) in (7.29) are \( O(1) \) as \( r \to 0 \) and the terms \( \Gamma \ast \Gamma \ast \Ric \) are \( O(r^2) \) as \( r \to 0 \) as needed. The other cases follow also by induction. \( \square \)

**Proof of Theorems 1.8 and 1.9.** By Lemma 7.6 we can find a change of coordinates \( \phi \) around the origin such that \( \phi_* g \) satisfies (7.17) and (7.18). Recall that the obstruction-flat systems have the form

\[
(7.31) \quad \Delta^{k-1} B_{ij} = \sum_{j=2}^{k+1} \sum_{\alpha_1 + \ldots + \alpha_j = 2(k+1)-2j} \nabla^\alpha_{g} Rm \ast \ldots \ast \nabla^\alpha_{g} Rm.
\]

From the expression (1.6) we can write the Bach tensor as

\[
(7.32) \quad B_{ij} = \Delta A_{ij} - \nabla_i \nabla^k A_{kj} + Rm \ast Rm,
\]

and using that \( g \) has constant scalar curvature together with the Bianchi identity we can rewrite (7.31) as

\[
(7.33) \quad \Delta^k g \Ric = \sum_{j=2}^{k+1} \sum_{\alpha_1 + \ldots + \alpha_j = 2(k+1)-2j} \nabla^\alpha_{g} Rm \ast \ldots \ast \nabla^\alpha_{g} Rm,
\]

which is exactly the same form as (1.13). From Lemma 7.1, (7.33) becomes

\[
(7.34) \quad \Delta^k g_{0 \ast \Ric} = T,
\]
where $T = O(r^{-2(k-1)})$ as $r \to 0$. Note that $T \in L^p$ near the origin for $p = \infty$ if $k = 1$ and for any $1 \leq p < \frac{n}{2(k-1)}$ if $k > 1$. On the other hand, since $\text{Ric}$ is bounded near the origin, $\text{Ric} \in L^p$ for any such $p$. It follows that $\text{Ric}(g)$ is a weak solution of $\Delta^k_{g_0} \text{Ric} = T$ on $B_\rho(0) \setminus \{0\}$ and it is easy to prove that in that case $\text{Ric}(g)$ extends to a weak solution of (7.33) on $B_\rho(0)$. We conclude that $\text{Ric} \in W^{2k,p}$. Choose

$$p = \begin{cases} 
\infty & \text{if } k = 1 \\
(1 - \epsilon)\frac{n}{2(k-1)} & \text{with } 0 < \epsilon < \frac{1}{2k-1} \text{ if } k > 1.
\end{cases}$$

(7.35)

Observe also that $0 < 2k - 1 - \frac{2}{p} \leq 1$ so by the Sobolev inequality, for any $\alpha$ such that $0 < \alpha < 2k - 1 - \frac{n}{p}$ we have

$$\|\nabla \text{Ric}\|_{C^\alpha(B_\rho(0))} \leq C\|\nabla \text{Ric}\|_{W^{2k-1,p}(B_\rho(0))},$$

(7.36)

with $C = C(n,p,k,\rho)$ and then $\text{Ric} \in C^{1,\alpha}(B_\rho(0))$. Note that from the estimates (7.17) and (7.18) we have $g \in C^{1,\alpha}$, which is sufficient for the existence of harmonic coordinates at the origin [DK81, Lemma 1.2]. In this harmonic coordinate system the metric $g$ is also $C^{1,\alpha}$ and solves (7.33). Moreover, by (7.36) and [DK81, Corollary 1.4], $\text{Ric} \in C^{1,\alpha}$ near the origin.

We then have that $g$ is a solution of (7.33) and is also a solution of an equation of the form

$$\frac{1}{2}g^{ij}\partial^2_{ij}g_{kl} + Q_{kl}(\partial g, g) = -\text{Ric}_{kl}(g),$$

(7.37)

where $Q(\partial g, g)$ is an expression that is quadratic in $\partial g$, polynomial in $g$ and has $\sqrt{|g|}$ in its denominator. Letting $p$ and $\alpha$ be as above, we know that $g$ and $\text{Ric}(g)$ are $C^{1,\alpha}$ at the origin, in particular they both are in $W^{1,p}$. Using elliptic regularity in (7.37) we conclude that $g \in W^{3,p}$ and the Sobolev inequality (compare (7.36)) implies that $g \in C^{2,\alpha}$. Furthermore, since $\text{Ric} \in C^{1,\alpha}$ it also follows that $g \in C^{3,\alpha}$ (see [DK81, Theorem 4.5]). With this regularity in $g$ we can write (7.33) as (7.34) in harmonic coordinates, i.e., we can write (7.33) as an equation of the form $\Delta^k_{g_0} \text{Ric} = T'$ with $T' \in L^p$.

Next, we claim that $\Delta^k_{g_0} \text{Ric} \in W^{1,p}$. To see this, take one derivative of (7.33). Since $g \in C^{3,\alpha}$, one sees that all of the terms on the right hand side are $O(r^{-2(k-1)})$ as $r \to 0$, which is in $L^p$ for $p$ as in (7.35). This shows that $\Delta^k_{g_0} \text{Ric} \in W^{1,p}$. Replacing $\Delta^k_{g} \text{Ric}$ with $\Delta^k_{g_0} \text{Ric}$ will introduce terms as in Lemma [7.7] but using the fact that $\text{Ric}$ is now in $C^{1,\alpha}$, we see that these terms are also in $W^{1,p}$, and consequently $\Delta^k_{g_0} \text{Ric} \in W^{1,p}$. Elliptic regularity then implies that $\text{Ric} \in W^{2k+1,p}$. By the Sobolev inequality, $\text{Ric} \in C^{2,\alpha}$ and then (7.37) implies $g \in C^{4,\alpha}$. It is clear that we can bootstrap the above argument to prove that $g$ is smooth at the origin. \qed
Appendix A. Turan’s Lemma

Given complex numbers \( z_1, \ldots, z_d \) and a nonnegative integer \( l \) we use \( S_l \) denote the sum

\[
c_1, \ldots, c_d \in \mathbb{C}, \quad S_l = \sum_{j=1}^{d} c_j z_j^l.
\]

The following lemma is the version of Turan’s Lemma that we will use throughout this section:

**Lemma A.1.** Let \( z_1, \ldots, z_d \) where \( d > 1 \) be complex numbers with \( |z_j| \geq 1 \) for \( j = 1, \ldots, d \) and let \( m \) be an integer with \( m \geq 1 \). Then

\[
|S_0|^2 \leq C \max\{|S_{m+1}|^2, \ldots, |S_{m+d}|^2\},
\]

where the positive constant \( C = C(m, d) \) can be estimated as

\[
C \leq A(d) \left( \frac{m + d}{d} \right)^{2(d-1)},
\]

for some constant \( A(d) \) depending only on \( d \).

**Proof.** The proof is a straightforward modification of that in [Naz93, Section 1.1]. □

A.1. An integral form of Turan’s Lemma. Let \( h \) be a positive number and consider the arithmetic progression \( \{t_l = lh\}_l \) where \( l \) are nonnegative integers. Consider also complex numbers \( \zeta_1, \ldots, \zeta_d \) such that \( \text{Re}(\zeta_j) \geq 0 \) for \( j = 1, \ldots, d \). If we let \( p(t) \) denote

\[
p(t) = \sum_{j=1}^{d} c_j e^{\zeta_j t},
\]

then

\[
p(lh) = \sum_{j=1}^{d} c_j z_j^l(h),
\]

where \( z_j(h) = e^{\zeta_j h} \). Since \( |z_j(h)| \geq 1 \) we have from Lemma A.1 an inequality of the form

\[
|p(0)|^2 \leq C(m, d) \left( \max_{l=m+1, \ldots, m+d} \{|p(lh)|^2\} \right) \leq C(m, d) \left( \sum_{l=m+1}^{m+d} |p(lh)|^2 \right),
\]

for any integer \( m \geq 1 \) where \( C(m, d) \leq A(d) \left( \frac{m+d}{d} \right)^{2(d-1)} \).

**Lemma A.2.** Let \( \zeta_1, \ldots, \zeta_d \) be complex numbers satisfying \( \text{Re}(\zeta_j) \geq 0 \) for \( j = 1, \ldots, d \), and let \( p(t) = \sum_{j=1}^{d} c_j e^{\zeta_j t} \). Then for any positive numbers \( 0 < a < b \)

\[
|p(0)|^2 \leq A(d) \left( \frac{b}{b-a} \right)^{2(d-1)} \frac{(b+a)}{(b-a)^2} \int_{a}^{b} |p(t)|^2 dt.
\]
Proof. Let \( c = \frac{a+b}{d} \), set \( h_0 = \frac{b-ca}{b-a} \) and let \( m \) be the integer part of \( \frac{c}{h_0} \) (i.e. the only integer \( m \) such that \( m \leq \frac{c}{h_0} < m+1 \)). Note that in this case \( m \geq 1 \). For any \( h > 0 \) by (A.6) we have

\[
(A.8) \quad |p(0)|^2 \leq A(d) \left( \frac{m+d}{d} \right)^{2(d-1)} \sum_{l=m+1}^{m+d} \left( |p(lh)|^2 \right).
\]

By our choice of \( m \) and \( h_0 \) we have

\[
(A.9) \quad (m+d)h_0 \leq c + dh_0 = b,
\]

so then \( a \leq lh_0 \leq b \) for \( l = m + 1, \ldots, m + d \). Let \( \eta = \frac{2a}{a+b} \), i.e. \( c\eta = a \). Note that \( 0 < \eta < 1 \). By taking the average of (A.8) respect to \( h \) on the interval \([\eta h_0, h_0]\) we have

\[
(A.10) \quad |p(0)|^2 \leq A(d) \left( \frac{m+d}{d} \right)^{2(d-1)} \frac{1}{h_0(1-\eta)} \int_{\eta h_0}^{h_0} \left( \sum_{l=m+1}^{m+d} |p(lh)|^2 \right) dh,
\]

It is easy to see that if \( f : \mathbb{R} \to \mathbb{R} \) is a continuous nonnegative function and and \( T_0, T_1 \) are positive numbers with \( T_0 < T_1 \), then

\[
(A.11) \quad \int_{T_0}^{T_1} \left( \sum_{l=m+1}^{m+d} f(lh) \right) dh \leq \frac{d}{2} \int_{(m+1)T_0}^{(m+d)T_1} f(h) dh.
\]

Therefore, we have

\[
(A.12) \quad |p(0)|^2 \leq A(d) \frac{d}{2} \left( \frac{m+d}{d} \right)^{2(d-1)} \frac{1}{h_0(1-\eta)} \int_0^{h_0} |p(t)|^2 dt,
\]

where we have used that \((m+1)\eta h_0 > c\eta = a\). We also have \((1-\eta)h_0 = \frac{(b-a)^2}{2d(b+a)}\). On the other hand

\[
(A.13) \quad \left( \frac{m+d}{d} \right)^{2(d-1)} = \left( \frac{(m+d)h_0}{dh_0} \right)^{2(d-1)} \leq 2^{2(d-1)} \left( \frac{b}{b-a} \right)^{2(d-1)},
\]

so the result follows.

Corollary A.3. If \( p(t) \) is as before then for any \( R > 0 \) we have

\[
(A.14) \quad \|p\|_{L^\infty[0,R]}^2 \leq \frac{A(d)}{R} \int_{2R}^{2R} |p(t)|^2 dt,
\]

and also

\[
(A.15) \quad \int_0^R |p(t)|^2 dt \leq A(d) \int_{3R}^{2R} \frac{2R}{\frac{3R}{2}} |p(t)|^2 dt.
\]

Proof. Let \( t_0 \in [0,R] \) be such that \( \|p\|_{L^\infty[0,R]} = |p(t_0)| \), and consider \( p_{t_0}(\tau) \) given by

\[
(A.16) \quad p_{t_0}(\tau) = \sum_{j=1}^d c_j e^{\xi_j t_0} e^{\xi_j \tau}.
\]
Clearly \( p_{t_0}(t - t_0) = p(t) \) and \( |p_{t_0}(0)| = \|p\|_{L^\infty[0,R]} \). By Lemma A.2 we have

\[
|p_{t_0}(0)|^2 \leq 4 \cdot 2^{d-1} A(d) \left( \frac{2R - t_0}{R} \right)^{d-1} \left( \frac{7}{2R} - 2\frac{t_0}{R^2} \right) \int_{4R - 2t_0}^{2R - t_0} |p_{t_0}(\tau)|^2 d\tau.
\]

(A.17)

Using a change of variables in (A.17) we obtain

\[
\|p\|^2_{L^\infty[0,R]} \leq 4 \cdot 4^{d-1} A(d) \left( \frac{7}{2R} \right) \int_{\frac{3}{2}R}^{2R} |p(t)|^2 dt,
\]

(A.18)

which proves (A.14). Once we have shown (A.14) we can prove (A.15) by writing

\[
\frac{1}{R} \int_0^R |p(t)|^2 dt \leq \|p\|^2_{L^\infty[0,R]} \leq \frac{A(d)}{R} \int_{\frac{3}{2}R}^{2R} |p(t)|^2 dt.
\]

(A.19)

\[\Box\]

A.2. Proof of Lemma 5.2. By making the change of variables \( r = e^t \), it is clear from (5.23) that it suffices to show the following lemma

Lemma A.4. Let \( p(t) \) be a sum of the form

\[
p(t) = \sum_{j=1}^d \sum_{s=0}^{n_j} c_{j,s} t^s e^{\zeta_j t},
\]

(A.20)

where \( n_j \geq 0 \) are integers, and \( c_{j,s} \in \mathbb{C} \) are fixed. Let \( M = \sum_{j=1}^d n_j \). Then

1. If \( \lambda = \min\{\text{Re}(\zeta_1), \ldots, \text{Re}(\zeta_d)\} > 0 \) then for every positive integer \( l \),

\[
e^{\lambda R} \int_{(l-1)R}^{lR} |p(t)|^2 dt \leq A(M + d) \int_{lR}^{(l+1)R} |p(t)|^2 dt.
\]

(A.21)

2. If \( \lambda = \min\{-\text{Re}(\zeta_1), \ldots, -\text{Re}(\zeta_d)\} > 0 \), then for any positive integer \( l \),

\[
\int_{lR}^{(l+1)R} |p(t)|^2 dt \leq C(M + d) e^{-\lambda R} \int_{(l-1)R}^{lR} |p(t)|^2 dt.
\]

(A.22)

Proof. We first consider functions \( p : \mathbb{R} \mapsto \mathbb{C} \) that have the form

\[
p(t) = \sum_{j=1}^d c_j e^{\zeta_j t},
\]

(A.23)

with \( c_j \in \mathbb{C} \), i.e., we do not consider the case of roots \( \zeta_j \) with multiplicities. Suppose that all numbers \( \zeta_j, j = 1, \ldots, d \) have positive real part. Let us first prove (A.21) for \( l = 1 \). Consider the function \( \tilde{p}(t) \) given by

\[
\tilde{p}(t) = \sum_{j=1}^d c_j e^{(\zeta_j - \lambda) t}.
\]

(A.24)
Since $Re(\zeta_j - \lambda) \geq 0$ for all $j = 1, \ldots, d$, we have from Corollary [A.3]

\begin{equation}
(A.25) \quad \int_0^R |\tilde{p}(t)|^2 dt \leq A(d) \int_{\frac{R}{2}}^{2R} |\tilde{p}(t)|^2 dt,
\end{equation}

multiplying both sides of (A.25) by $e^{3\lambda R}$ we have

\begin{equation}
(A.26) \quad e^{\lambda R} \int_0^R e^{2\lambda t} |\tilde{p}(t)|^2 dt \leq e^{3\lambda R} \int_0^R |\tilde{p}(t)|^2 dt \leq A(d) e^{3\lambda R} \int_{\frac{R}{2}}^{\frac{R}{2}} |\tilde{p}(t)|^2 dt \leq A(d) \int_R^{2R} e^{2\lambda t} |\tilde{p}(t)|^2 dt,
\end{equation}

and $e^{2\lambda t} |\tilde{p}(t)|^2 = |p(t)|^2$. For the case $l > 1$, we write any $t \in [(l - 1)R, (l + 1)R]$ as $t = (l - 1)R + \tau$ where $\tau \in [0, 2R]$ and then we write $p(t)$ as

\begin{equation}
(A.27) \quad p(t) = q(\tau) = \sum_{j=1}^d c_j e^{\zeta_j \tau},
\end{equation}

where $c_{j,l} = c_j e^{(l-1)R}$. Applying the above argument to $q(\tau)$ then (A.21) follows after a change of variables. If now all numbers $Re(\zeta_j)$ are negative for $j = 1, \ldots, d$, it suffices to prove (A.22) for $l = 1$ since as before, the general case $l \geq 1$ follows after a change of variables. Let $\tilde{p}(t) = \sum_{j=1}^d c_j e^{(\zeta_j + \lambda) t}$ and write $t \in [0, 2R]$ as $t = 2R - \tau$ where $\tau \in [0, 2R]$, then $c_j e^{(\zeta_j + \lambda) t} = c_j e^{(\zeta_j + \lambda) 2R} e^{-(\zeta_j + \lambda) \tau}$ and $Re(- (\zeta_j + \lambda)) \geq 0$, so by (A.21), if we let $q(\tau) = \sum_{j=1}^d c_j e^{(\zeta_j + \lambda) 2R} e^{-(\zeta_j + \lambda) \tau}$ we obtain

\begin{equation}
(A.28) \quad \int_0^R |q(\tau)|^2 d\tau \leq A(d) \int_{\frac{R}{2}}^{2R} |q(\tau)|^2 d\tau.
\end{equation}

On the other hand,

\begin{equation}
(A.29) \quad \int_0^R |q(\tau)|^2 d\tau = \int_R^{2R} |\tilde{p}(t)|^2 dt, \quad \text{and} \quad \int_{\frac{R}{2}}^{2R} |q(\tau)|^2 d\tau = \int_0^{\frac{R}{2}} |\tilde{p}(t)|^2 dt,
\end{equation}

so from (A.28), we have

\begin{equation}
(A.30) \quad e^{\lambda R} \int_R^{2R} e^{-2\lambda t} |\tilde{p}(t)|^2 dt \leq e^{-\lambda R} \int_R^{2R} |\tilde{p}(t)|^2 dt \leq A(d) e^{-\lambda R} \int_{\frac{R}{2}}^{\frac{R}{2}} |\tilde{p}(t)|^2 dt \leq A(d) \int_{\frac{R}{2}}^{\frac{R}{2}} e^{-2\lambda t} |\tilde{p}(t)|^2 dt \leq A(d) \int_0^R e^{-2\lambda t} |\tilde{p}(t)|^2 dt,
\end{equation}

and $e^{-2\lambda t} |\tilde{p}(t)|^2 = |p(t)|^2$.

For the general case involving multiplicity, we will only prove the statement for the case $\min\{Re(\zeta_1), \ldots, Re(\zeta_d)\} > 0$ since the other case will follow as in the above argument. We consider first the case corresponding to $n_1 = 1$ and $n_2 = \ldots = n_d = 0$. 


Let $\epsilon > 0$ and let $p_{\epsilon}(t)$ be given by

\begin{align}
(A.31) \quad p_{\epsilon}(t) &= c_{1,0}e^{\xi_{1}t} + \frac{c_{1,1}}{\epsilon}(e^{(\xi_{1}+\epsilon)t} - e^{\xi_{1}t}) + \sum_{j=2}^{d} c_{j,0}e^{\xi_{j}t},
\end{align}

then by (A.21) for the case with no multiplicity, we have

\begin{align}
(A.32) \quad e^{\lambda R} \int_{(l-1)R}^{lR} |p_{\epsilon}(t)|^{2} dt \leq A(d+1) \int_{lR}^{(l+1)R} |p_{\epsilon}(t)|^{2} dt,
\end{align}

where $\lambda$ is defined as before, and since $A(d+1)$ does not depend on $\epsilon$ we can take the limit of (A.32) as $\epsilon$ tends to zero and obtain (A.21). For the higher multiplicity case, (A.21) can be proved using induction. \hfill \Box

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