INITIAL DATA ON BIG BANG SINGULARITIES

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Abstract. The goal of this article is to parametrise solutions to Einstein’s equations with big bang singularities and quiescent asymptotics, disregarding issues such as spikes. To this end, we introduce a notion of initial data on big bang singularities and conjecture that it can be used to parametrise quiescent solutions. A mathematical statement of the conjecture presupposes a precise definition of the class of quiescent solutions as well as a proof of existence and uniqueness of developments corresponding to initial data on a big bang singularity. Needless to say, there many possible definitions. Nevertheless, we provide one definition of quiescence which is satisfied by many examples and that leads to initial data of the type considered here. We also appeal to existing results in order to illustrate that, in certain cases, there are unique developments corresponding to initial data on the singularity. However, the perspective developed here naturally leads to a large class of open problems corresponding to the general conjecture. Moreover, there are many versions of the conjecture in that one can consider different matter models as well as different foliations of the spacetimes in the vicinity of a big bang singularity. An additional benefit of the notion of initial data developed here is that it can be used to give a unified perspective on the existing results concerning quiescent singularities. In fact, we provide several examples of how existing results can be considered to be special cases of the framework developed here. A second, potential, application is to oscillatory and spatially inhomogeneous big bang singularities. Considering the existing arguments in the spatially homogeneous setting, a crucial first step in the study of oscillatory behaviour is to understand how solutions approach the Kasner circle along a stable manifold, and then depart via an unstable manifold. In order to carry out a similar analysis in the spatially inhomogeneous setting, it is of central importance to first identify the stable manifold. Building on the work of Fournodavlos and Luk, we here propose such an identification.

1. Introduction

In the last 20 years or so, a large number of results have appeared in which the authors specify data on a big bang singularity and then prove the existence of corresponding solutions; cf., e.g., [44, 41, 57, 6, 43, 75, 1, 46, 2, 29]. The solutions are typically crushing (in the sense that with respect to a suitable foliation, the mean curvature, say \( \theta \), diverges uniformly to infinity in the direction of the singularity) and quiescent (in the sense that if \( \bar{K} \) is the Weingarten map of the leaves of the foliation, then the expansion normalised Weingarten map, defined by \( K := \bar{K}/\theta \), has the property that its eigenvalues converge in the direction of the singularity). However, what the data on the singularity are depends on the situation under consideration. In some cases, the data correspond to the asymptotics of a symmetry reduced version of the equations, in some cases the initial data correspond to solutions to a simplified version of the equations etc. Moreover, there is no unified geometric perspective explaining when quiescent asymptotics are to be expected, and thereby also when they are not to be expected. Finally, in order to understand the oscillatory setting (roughly speaking, the case that the eigenvalues of \( K \) do not converge), it is of interest to identify the stable manifold associated with the Kasner circle for Einstein’s vacuum equations in 3 + 1-dimensions, in the absence of symmetries; cf. Subsection 2.2 below for a more detailed discussion of this aspect. The purpose of this article is to address all of these issues, partly on a conjectural level, and the notion of initial data on a singularity we introduce here is the main tool.

1.1. Quiescent solutions to Einstein’s vacuum equations without symmetries. The notion of initial data on a singularity we propose here arises naturally when considering the recent results of Fournodavlos and Luk, [29], with the geometric perspective developed in [66] in mind.
For this reason, it is natural to begin by a discussion of the results of [29]. In this article, the authors specify the following initial data on the singularity

**Definition 1.** Fournodavlos and Luk initial data on the singularity consist of smooth functions $c_{ij}, p_i : T^3 \to \mathbb{R}$, $i,j = 1,2,3$ satisfying the following conditions: $c_{ij} = c_{ji}$; $p_1 < p_2 < p_3$; $c_{11}, c_{22}, c_{33} > 0$;

(1) $\sum_{i=1}^{3} p_i = \sum_{i=1}^{3} p_i^2 = 1$;

and

(2) $\sum_{i=1}^{3} \left[ \partial_i c_{il}(p_i - p_i) + 2 \partial_i \kappa^l_i + \mathbb{I}_{l>i} \frac{\partial_j (c_{11} c_{22} c_{33})}{c_{11} c_{22} c_{33}} \kappa^j_i \right] = 0$

for $i = 1,2,3$. Here $\kappa^j_i = -p_i$ (no summation); $\kappa^j_i = 0$ if $l < i$;

(3) $\kappa_1^2 := p_1 - 1$; $\kappa_2^3 := (p_2 - p_3) \frac{c_{12}^2}{c_{22}}$; $\kappa_3^1 := (p_3 - p_1) \frac{c_{12}^2}{c_{22} c_{33}} + (p_3 - p_1) \frac{c_{13}^2}{c_{33}}$;

$\mathbb{I}_{l>i} = 1$ if $l > i$; and $\mathbb{I}_{l>i} = 0$ if $l \leq i$.

Given data satisfying these conditions, there is, due to [29] Theorem 1.1, p. 2, a $C^2$ solution to the Einstein vacuum equations $\text{Ric} = 0$ of the form

(4) $g = -dt \otimes dt + \sum_{i,j=1}^{3} a_{ij} t^{2 \max(i,j)} dx^i \otimes dx^j$,

where $(t, x^1, x^2, x^3) \in (0, T) \times T^3$ for some $T > 0$ and $a_{ij} : (0, T) \times T^3 \to \mathbb{R}$ are smooth functions (symmetric in $i$ and $j$) which extend to continuous functions : $[0, T] \times T^3 \to \mathbb{R}$. Moreover, the $a_{ij}$ obey

(5) $\lim_{t \to 0^+} a_{ij}(t, x) = c_{ij}(x)$.

However, in the course of the argument, the authors also obtain the following control of the Weingarten map:

(6) $\sum_{\tau = 0}^{1} \sum_{|\kappa| \leq 2-\tau} |t^{\tau+1} |\partial_\tau \partial_\tau (k_j^i - t^{-1} \kappa^j_i)| = O(\text{min}\{t^\epsilon, t^{-2p_i+2p_j}\})$,

where $\epsilon > 0$ and $k$ is the second fundamental form of the leaves of the foliation (using the conventions of [29]; here we use the opposite conventions). Clearly, the formulation of this result is tied quite closely to the specific situation under consideration: $T^3$ spatial topology; Gaussian foliation etc. On the other hand, the results of [65, 66] indicate that quiescent behaviour in the case of 3 + 1-dimensional vacuum solutions to Einstein’s equations should be tied to properties of the expansion normalised Weingarten map. It is therefore of interest to try to combine the formulation in [29] with the more geometric perspective taken in [66].

### 1.2. A geometric formulation of initial data on the singularity

In order to relate the results of [29] with those of [65, 66], note that our conventions concerning the second fundamental form are opposite to those of [29]. Here we use $k$ to denote the second fundamental form of the leaves of the foliation, and $K = -k$, where $k$ is object appearing in [6]. Moreover, we denote the Weingarten map by $\hat{K}$; i.e., $\hat{K}$ is the $(1,1)$-tensor field satisfying $\hat{K}^j_i = k_j^i$. If $\theta$ denotes the mean curvature of the leaves of the foliation, the expansion normalised Weingarten map is defined by $\hat{\mathcal{K}} := \hat{K}/\theta$, assuming $\theta \neq 0$. Note here that taking the trace of (6), keeping the definition of $\kappa$ and the conventions in mind, yields $t\theta - 1 = O(t^\epsilon)$ in $C^2$. Combining this observation with [6] again yields the conclusion that $\hat{\mathcal{K}}$ converges to a limit, say $\mathcal{K}$, in $C^2$, and the rate of convergence is $t^\epsilon$. In fact, $\mathcal{K}^j_i = -k_j^i$. At this stage, it is of interest to reformulate the conditions of Definition 1 in terms of $\mathcal{K}$. Note, to this end, that (1) can be reformulated to $\text{tr.}\mathcal{K} = \text{tr.}\mathcal{K}^2 = 1$. The fact that the $p_i$ are distinct can be reformulated as saying that the eigenvalues of $\mathcal{K}$ are distinct. Next, note that (4) and (5) are closely tied to the coordinate vector fields on $T^3$. From a geometric
Combining the above terminology with (3) and (5), it can also be verified that

$$
\text{(7a)} \quad \mathcal{F}_3 = \partial_3, \quad \mathcal{F}_2 = \partial_2 + \mathcal{F}_2^2 \partial_3, \quad \mathcal{F}_1 = \partial_1 + \mathcal{F}_1^2 \partial_2 + \mathcal{F}_1^3 \partial_3,
$$

$$
\text{(7b)} \quad \mathcal{F}^1 = dx^1, \quad \mathcal{F}^2 = dx^2 - \mathcal{F}_1^2 dx^1, \quad \mathcal{F}^3 = dx^3 - \mathcal{F}_2^3 dx^2 + (\mathcal{F}_1^2 \mathcal{F}_2^3 - \mathcal{F}_1^3 \mathcal{F}_2^2) dx^1,
$$

where

$$
\text{(8)} \quad \mathcal{F}_1^2 = \frac{1}{p_1 - p_2} \mathcal{X}_1^2, \quad \mathcal{F}_2^3 = \frac{1}{p_2 - p_3} \mathcal{X}_2^3, \quad \mathcal{F}_3^3 = \frac{1}{p_1 - p_3} \left( \mathcal{X}_1^3 + \frac{1}{p_1 - p_2} \mathcal{X}_1^1 \mathcal{X}_2^2 \right).
$$

With this notation, \( \mathcal{K} \mathcal{F}_A = p_A \mathcal{F}_A \) (no summation). Moreover, (4) can be written

$$
\text{(9)} \quad g = -dt \otimes dt + \sum_{A,B=1}^3 b_{AB} t^2 p_{\max(A,B)} \mathcal{F}^A \otimes \mathcal{F}^B.
$$

Combining the above terminology with (3) and (5), it can also be verified that

$$
\lim_{t \to 0^+} b_{AB} = c_{AA} \delta_{AB}
$$

(no summation on \(A\)), so that

$$
\text{(10)} \quad \lim_{t \to 0^+} \sum_{A,B=1}^3 b_{AB} \mathcal{F}^A \otimes \mathcal{F}^B = \sum_{A=1}^3 c_{AA} \mathcal{F}^A \otimes \mathcal{F}^A =: \bar{h}.
$$

Since \(c_{AA} > 0\), the limit, i.e. \( \bar{h} \), is a Riemannian metric on \( \mathbb{T}^3 \). At this stage, we thus have two natural structures on \( \mathbb{T}^3 \): the \((1,1)\)-tensor field \( \mathcal{K} \) and the Riemannian metric \( \bar{h} \). It would of course be desirable if all the conditions could be expressed in terms of these quantities. From the above observations it is clear that the \( p_i \) can be read off as the eigenvalues of \( \mathcal{K} \); that \( (1) \) can be reformulated to \( \text{tr} \mathcal{K} = \text{tr} \mathcal{K}_2 = 1 \); and that \( \mathcal{K} \) is symmetric with respect to \( \bar{h} \), cf. (10). Next, on the basis of (6), the structure “constants” \( \gamma^A_{BC} \) associated with the frame \( \{ \mathcal{F}_A \} \) can be expected to be important. They are defined by

$$
\text{(11)} \quad [\mathcal{F}_B, \mathcal{F}_C] = \gamma^A_{BC} \mathcal{F}_A.
$$

Due to (6), we would expect \( \gamma^A_{23} = 0 \) to be a necessary condition in order to obtain the quiescent behaviour of (29) (note that this condition is independent of the normalisation of the eigenvector fields and that it is crucially dependent on the assumption that \( p_1 < p_2 < p_3 \)). However, the frame introduced in (7) satisfies this condition. In fact, this condition is automatically built into the definition of the initial data. Finally, we need to express (2) geometrically. Since this equation is the limit of the momentum constraint, cf. (29) Remark 1.4, p. 3], it is of interest to relate it to \( \text{div}_h \mathcal{K} \). In fact, it turns out that (2) is equivalent to the condition that \( \text{div}_h \mathcal{K} = 0 \); cf. Lemma (29) below. The above observations suggest the following definition of initial data on a big bang singularity for Einstein’s vacuum equations in \( 3 + 1 \)-dimensions.

**Definition 2.** Let \((\bar{M}, \bar{h})\) be a smooth 3-dimensional Riemannian manifold and \( \mathcal{K} \) be a smooth \((1,1)\)-tensor field on \( \bar{M} \). Then \((\bar{M}, \bar{h}, \mathcal{K})\) are non-degenerate quiescent vacuum initial data on the singularity if the following holds:

1. \( \text{tr} \mathcal{K} = 1 \) and \( \mathcal{K} \) is symmetric with respect to \( \bar{h} \).
2. \( \text{tr} \mathcal{K}_2 = 1 \) and \( \text{div}_h \mathcal{K} = 0 \).
3. The eigenvalues of \( \mathcal{K} \) are distinct.
4. \( \mathcal{O}^1_{23} \equiv 0 \) on \( \bar{M} \), where (if \( p_1 < p_2 < p_3 \) are the eigenvalues of \( \mathcal{K} \), \( \mathcal{F}_A \) is an eigenvector field corresponding to \( p_A \), normalised so that \( \mathcal{F}_A|_p = 1 \), and \( \gamma^A_{BC} \), is defined by \( [\mathcal{F}_B, \mathcal{F}_C] = \gamma^A_{BC} \mathcal{F}_A \) the quantity \( \mathcal{O}^A_{BC} \) is defined by \( \mathcal{O}^A_{BC} := (\gamma^A_{BC})^2 \).

**Remark 3.** That \( \mathcal{K} \) is symmetric with respect to \( \bar{h} \) means that if \( p \in \bar{M} \) and \( \xi, \zeta \in T_p \bar{M} \), then \( \bar{h}(\xi, \zeta) = \bar{h}(\xi, \mathcal{K} \zeta) \).

**Remark 4.** Due to the normalisation of the eigenvector fields, they are uniquely defined up to a sign. For this reason, the functions \( \mathcal{O}^A_{BC} \) are uniquely defined smooth functions on \( \bar{M} \).

Due to the observations made prior to the statement of the definition, the following holds.
Proposition 5. Given Fournodavlos and Luk initial data on the singularity in the sense of Definition 1, define $\mathcal{K}$ by $\mathcal{K}^i_j = -\kappa^i_j$ and $\hat{h}$ by (I(d)). Then $(\mathbb{T}^3, \hat{h}, \mathcal{K})$ are non-degenerate quiescent vacuum initial data on the singularity.

In short, initial data in the sense of Definition 1 are a special case of initial data in the sense of Definition 2. It is therefore of interest to find out if it is possible to go in the other direction. On the level of local coordinates, it turns out to be possible.

Proposition 6. Let $(M, \bar{h}, \mathcal{K})$ be non-degenerate quiescent vacuum initial data on the singularity in the sense of Definition 2 and let $p \in M$. Let $p_A$, $A = 1, 2, 3$, be the eigenvalues of $\mathcal{K}$, and assume them to be ordered so that $p_1 < p_2 < p_3$. Then there are local coordinates $(V, \bar{x})$ such that $p \in V$ and such that $\mathcal{K}$ has a frame $\{ \mathcal{X}_A \}$ of eigenvector fields on $V$ given by (7a) and (8), where the sub- and superscripts refer to the coordinates $(V, \bar{x})$. Define $c_{AA} := \hat{h}(\mathcal{X}_A, \mathcal{X}_A)$ (no summation). Define

\[
\begin{align*}
    c_{12} := & \frac{c_{33}}{p_2 - p_1} \mathcal{K}^1_2, \\
    c_{23} := & \frac{c_{33}}{p_3 - p_2} \mathcal{K}^2_3, \\
    c_{13} := & \frac{c_{33}}{p_3 - p_1} \left( \mathcal{K}^1_3 + \frac{1}{p_3 - p_2} \mathcal{K}^1_2 \mathcal{K}^2_3 \right),
\end{align*}
\]

where the components of $\mathcal{K}$ are calculated with respect to the local coordinates $(V, \bar{x})$, and define the remaining components $c_{ij}$ by requiring $c_{ij} = c_{ji}$. Then the functions $p_i$ and $c_{ij}$ satisfy all the conditions of Definition 1, but on the set $V$ instead of $\mathbb{T}^3$, and where the sub-, superscripts and derivatives refer to the coordinate frame associated with $(V, \bar{x})$.

Proof. The proof of the statement is to be found at the end of Subsection 3.2 \hfill $\square$

This observation indicates that Definitions 1 and 2 are locally the same. On the other hand, it is conceivable that a more general definition of data on the singularity is possible and would, non-trivially, include both Definitions 1 and 2. We return to this topic in Subsection 1.7 below.

One potential deficiency in the definitions is the non-degeneracy condition. In some situations, it is to be expected that this condition can be relaxed; cf. Remark 26 below. On the other hand, one of the main applications of the definitions is to think of the data on the singularity as parametrising the stable manifold associated with the Kasner circle. One of the main potential uses of this perspective is to the understanding of oscillatory spatially inhomogeneous spacetimes; cf. Subsection 2.2 below. In that context, the degenerate case is not of interest; cf. Subsection 2.2 below.

1.3. Developments of non-degenerate quiescent vacuum initial data on the singularity.

Once one has specified initial data on the singularity, the question arises how these data are related to a corresponding development. In the case of [29], we have already given the relation; cf. (4)–(6). Moreover, [9, 10] and the fact that $\mathcal{K} \to \mathcal{K}$ suggest how to relate the development with the data more generally. In the setting of Definition 2 we use the following terminology.

Definition 7. Let $(M, \bar{h}, \mathcal{K})$ be non-degenerate quiescent vacuum initial data on the singularity in the sense of Definition 2. A locally Gaussian development of the initial data is then a time oriented Lorentz manifold $(M, g)$, solving Einstein’s vacuum equations, such that the following holds. There is a $0 < t_+ \in \mathbb{R}$ and a diffeomorphism $\Psi$ from $M \times (0, t_+)$ to an open subset of $(M, g)$ such that

\[
\Psi^* g = -dt \otimes dt + \sum_{A,B} b_{A,B} t^{2 \max \{A, B\}} \mathcal{X}^A \otimes \mathcal{X}^B,
\]

where $\{ \mathcal{X}_A \}$ is a local basis of eigenvector fields of $\mathcal{K}$, $\{ \mathcal{Y}^A \}$ is the dual basis, $\mathcal{K} \mathcal{X}_A = p_A \mathcal{X}_A$ (no summation) and $p_1 < p_2 < p_3$. Here $\Psi(M_t)$, where $M_t := M \times \{ t \}$, is required to have strictly positive mean curvature in $(M, g)$ for $t \in (0, t_+)$. Let $\theta$ be the mean curvature of the leaves of the foliation; $K$ be the Weingarten map of the leaves of the foliation; $K := K/\theta$ be the expansion normalised Weingarten map; and $\bar{h}$ be defined by

\[
\bar{h} := \sum_{A,B} b_{A,B} \mathcal{Y}^A \otimes \mathcal{Y}^B.
\]
note that \( \tilde{h} \) is globally well defined independent of the choice of local frame. Then the following correspondence between the solution and the asymptotic data is required to hold. There is an \( \varepsilon > 0 \) and for every \( 0 \leq t \in \mathbb{N} \), a constant \( C_t \) such that
\[
\left\| \mathcal{K}(\cdot, t) - \mathcal{X} \right\|_{C^1(M)} \leq C_t \varepsilon,
\]
\[
\left\| \tilde{h}(\cdot, t) - \tilde{h} \right\|_{C^1(M)} \leq C_t \varepsilon
\]
on \( M_+ := \tilde{M} \times (0, t_+) \). If, in addition to the above, \((M, g)\) is globally hyperbolic and the hypersurfaces \( \Psi(M_t) \) are Cauchy hypersurfaces in \((M, g)\) for \( t \in (0, t_+) \), then \((M, g)\) is called a locally Gaussian globally hyperbolic development of the initial data.

**Remark 8.** The convergence conditions (15) could both be weakened and strengthened; one could change the rate, change the norm, only insist on control of a finite number of derivatives etc. However, we here stick to this particular choice and encourage the reader to consider other possibilities.

**Remark 9.** Given initial data as in Definition [1], we locally obtain initial data in the sense of Definition [2] cf. Proposition [3]. Assuming it is possible to appeal to a localised version of [29], it is of interest to know if a development in the sense of [29] yields a development in the sense of Definition [7] (at least locally). That the answer to this question is yes is demonstrated in Proposition [61] below.

Comparing this definition with [29], it is of interest to note that the authors of [29] require control of the normal derivative of the second fundamental form in order to conclude uniqueness; cf. [29] Theorem 1.7, p. 4], in particular [29] (1.10), p. 4] (this estimate coincides with (1) above). They also require control of the asymptotic behaviour of \( \theta \). It is not immediately obvious how to estimate \( \theta \) by appealing to (15). However, it turns out to be possible to use these estimates to deduce that \( \theta - 1 \) converges to zero as \( t^e \) in \( C^i \) etc. We discuss this topic in greater detail in Subsection [1, 6] below.

1.4. **Generalisations.** So far, we have focused on the 3+1-dimensional vacuum setting. However, it is of interest to generalise Definitions [2] and [7]. For this reason, we now define data on the singularity in the case of higher dimensions, for scalar fields, and in the presence of a cosmological constant.

**Definition 10.** Let \( 3 \leq n \in \mathbb{N}, \Lambda \in \mathbb{R} \), \((\tilde{M}, \tilde{h})\) be a smooth \( n \)-dimensional Riemannian manifold, \( \mathcal{X} \) be a smooth \((1,1)\)-tensor field on \( \tilde{M} \) and \( \Phi_a \) be smooth functions on \( \tilde{M} \). Then \((\tilde{M}, \tilde{h}, \mathcal{X}, \Phi_a, \Phi_b)\) are non-degenerate quiescent initial data on the singularity for the Einstein-scalar field equations with a cosmological constant \( \Lambda \) if

1. \( \text{tr} \mathcal{X} = 1 \) and \( \mathcal{X} \) is symmetric with respect to \( \tilde{h} \).
2. \( \text{tr} \mathcal{X}^2 + \Phi_a^2 = 1 \) and \( \text{div} b \mathcal{X} = \Phi_a d \Phi_b \).
3. The eigenvalues of \( \mathcal{X} \) are distinct.
4. \( \mathcal{O}_{BC} \) vanishes in a neighbourhood of \( \tilde{x} \in \tilde{M} \) if \( 1 + p_A(\tilde{x}) - p_B(\tilde{x}) - p_C(\tilde{x}) \leq 0 \).

**Remark 11.** Here the \( \mathcal{O}_{BC} \) are defined as follows. Let \( p_1 < \cdots < p_n \) be the eigenvalues of \( \mathcal{X} \), let \( \mathcal{X}_A \) be an eigenvector field corresponding to \( p_A \) such that \( |\mathcal{X}_A|_{\tilde{h}} = 1 \), and define \( \mathbf{\gamma}_{BC}^A \) by \( \mathcal{O}_{BC} = \mathbf{\gamma}_{BC}^A \mathbf{\mathcal{X}}_A \). Then \( \mathcal{O}_{BC}^2 := (\mathbf{\gamma}_{BC}^A)^2 \).

**Remark 12.** Since \( \text{tr} \mathcal{X}^2 + \Phi_a^2 = 1 \), the sum of the \( p_A^2 \) is less than or equal to 1. In case \( n \geq 3 \), it is thus clear that \( |p_A| < 1 \) for all \( A \in \{1, \ldots, d\} \). In particular, \( 1 + p_A - p_B - p_C > 0 \) if \( A = B \) or if \( A = C \). In the last criterion of the definition, we can therefore assume that \( A \notin \{B, C\} \). Note also that \( \mathcal{O}_{BC}^2 = 0 \) if \( B = C \). To conclude, \( A, B \) and \( C \) can be assumed to be distinct in the last criterion of the definition. In case \( n = 3 \) and \( A, B \) and \( C \) are distinct, \( 1 + p_A - p_B - p_C = 2p_A \).

**Remark 13.** Condition 4 should be compared with the requirement \( 1 + p_1 - p_{n-1} - p_n > 0 \) appearing in [22] Theorem 10.1, p. 1104] (assuming \( p_1 < \cdots < p_n \)); cf. also [38] for the origin of this condition.
In case \( n = 3 \), the fourth condition is only a restriction if \( p_1(\hat{x}) \leq 0 \). Moreover, it then corresponds to the requirement that \( \partial^2_{t \hat{x}} = 0 \) in a neighbourhood of \( \hat{x} \). The reason for this is that if \( p_1 > 0 \), then \( p_i > 0 \), \( i = 1, 2, 3 \), and the fourth condition is void. If \( p_1 \leq 0 \), then \( p_2 > 0 \), since the sum of the \( p_i \) equal 1 and \( p_3 < 1 \). Thus the \( A \) appearing in the fourth condition has to be 1 and \( \{ B, C \} \) has to equal \( \{ 2, 3 \} \). This observation should be compared with condition 4 of Definition 2.

Remark 15. For \( A, B \) and \( C \) distinct, the condition that \( \partial^A_{BC} \) vanish is independent of the choice of normalisation of the eigenvector fields \( \mathcal{X}_A \).

Next, we generalise Definition 7.

Definition 16. Let \((\hat{M}, \hat{h}, \mathcal{X}, \Phi_\alpha, \Phi_\theta)\) be non-degenerate quiescent initial data on the singularity for the Einstein-scalar field equations with a cosmological constant \( \Lambda \), in the sense of Definition 10. A locally Gaussian development of the initial data is then a smooth time oriented Lorentz manifold \((\hat{M}, g)\) and a smooth function \( \phi \), solving the Einstein scalar field equations with a cosmological constant \( \Lambda \), such that the following holds. There is a \( 0 \leq t_+ \in \mathbb{R} \) and a diffeomorphism \( \Psi \) from \( \hat{M} \times (0, t_+) \) to an open subset of \((M, g)\) such that \( \Psi^* g \) can be written as in (13), where \( \{ \mathcal{X}_A \} \) and \( p_A \) are given in Remark 11 and \( \{ \Phi^A \} \) is the dual basis to \( \{ \mathcal{X}_A \} \). Moreover, there is a \( t_0 \in (0, t_+) \) and a \( 0 < \theta_0 \in \mathbb{R} \) such that \( \Psi(M_t) \), where \( M_t := \hat{M} \times \{ t \} \), has mean curvature \( \theta \) bounded from below by \( \theta_0 \) for \( t \leq t_0 \). If \( \Lambda > 0 \), \( \theta_0 \) is additionally required to satisfy \( \theta_0 > |2\Lambda/(n-1)|^{1/2} \). Let \( \mathcal{K} \) be the expansion normalised Weingarten map of the leaves of the foliation and \( \hat{h} \) be defined by (14). Then the following correspondence between the solution and the asymptotic data is required to hold. There is an \( \varepsilon > 0 \) and for every \( l \in \mathbb{N} \), a constant \( C_l \) such that

\[
\begin{align*}
(16a) & \quad \| \mathcal{K}(\cdot, t) - \mathcal{X} \|_{C^l(\hat{M})} \leq C_l t^\varepsilon, \\
(16b) & \quad \| \hat{h}(\cdot, t) - \hat{h} \|_{C^l(\hat{M})} \leq C_l t^\varepsilon, \\
(16c) & \quad \| (\hat{U} \phi)(\cdot, t) - \Phi_\alpha \|_{C^l(\hat{M})} \leq C_l t^\varepsilon, \\
(16d) & \quad \| (\phi - \Phi_\alpha \phi)(\cdot, t) - \Phi_\theta \|_{C^{l+1}(\hat{M})} \leq C_l t^\varepsilon
\end{align*}
\]

for all \( t \leq t_0 \). Here \( \hat{U} := \theta^{-1} \partial_t \). Moreover, \( \varrho \) denotes the logarithmic volume density, defined by the condition \( \varrho \varepsilon \mu_\hat{h} = \mu_\hat{g} \), where \( \mu_\hat{h} \) is the volume form on \( \hat{M} \) induced by \( \hat{h} \), \( \mu_\hat{g} \) is the volume form on \( \hat{M} \) induced by the metric \( \Psi^*\hat{g} \) on the leaves \( M_t \). If, in addition to the above, \((M, g)\) is globally hyperbolic and the hypersurfaces \( \Psi(M_t) \) are Cauchy hypersurfaces in \((M, g)\) for \( t \in (0, t_+) \), then \((M, g)\) is called a locally Gaussian globally hyperbolic development of the initial data.

It is sometimes of interest to consider foliations that are not Gaussian. Note, to this end, that given a development as in Definition 16 the metric can, in a neighbourhood of the singularity, be given a representation expressed in terms of the mean curvature instead of \( t \); cf. Theorem 35 in particular (29b), below. However, it might also be of interest to allow a lapse function different from 1. We therefore give a generalised definition of development, though we still require the shift vector field to vanish.

Definition 17. Let \((\hat{M}, \hat{h}, \mathcal{X}, \Phi_\alpha, \Phi_\theta)\) be non-degenerate quiescent initial data on the singularity for the Einstein-scalar field equations with a cosmological constant \( \Lambda \), in the sense of Definition 10. A local crushing development of the initial data with vanishing shift vector field is then a smooth time oriented Lorentz manifold \((\hat{M}, g)\) and a smooth function \( \phi \), solving the Einstein scalar field equations with a cosmological constant \( \Lambda \), such that the following holds. There is a \( 0 < t_+ \in \mathbb{R} \) and a diffeomorphism \( \Psi \) from \( \hat{M} \times (0, t_+) \) to an open subset of \((M, g)\) such that

\[
\Psi^* g = -N^2 dt \otimes dt + \sum_{A,B} h_{AB} \theta^{-2p_{\text{max}}(A,B)} \mathcal{X}^A \otimes \mathcal{X}^B,
\]

where \( N \) is a strictly positive function (the lapse function), \( \{ \mathcal{X}_A \} \) and \( p_A \) are given in Remark 11 and \( \{ \Phi^A \} \) is the dual basis to \( \{ \mathcal{X}_A \} \). Here \( \Psi(M_t) \), where \( M_t := \hat{M} \times \{ t \} \), is required to have strictly positive mean curvature \( \theta \) in \((M, g)\) for \( t \in (0, t_+) \). Moreover, \( \theta \) is required to diverge uniformly to \( \infty \) as \( t \to 0^+ \). Let \( \mathcal{K} \) be the expansion normalised Weingarten map of the leaves of the foliation and \( \hat{h} \) be defined by (14). Then the following correspondence between the solution
and the asymptotic data is required to hold. Then the following correspondence between the solution and the asymptotic data is required to hold. There is an \( \varepsilon > 0 \) and for every \( l \in \mathbb{N} \), a constant \( C_l \) such that

\[
(18a) \quad |\bar{D}^l (\mathcal{K} - \mathcal{K}^\prime)|_{\bar{h}} \leq C_l \theta^{-\varepsilon},
\]
\[
(18b) \quad |\bar{D}^l (\bar{h} - \bar{h}^\prime)|_{\bar{h}} \leq C_l \theta^{-\varepsilon},
\]
\[
(18c) \quad |\bar{D}^l (\bar{U} \phi - \Phi_0)|_{\bar{h}} \leq C_l \theta^{-\varepsilon},
\]
\[
(18d) \quad |\bar{D}^l (\phi - \Phi_0 q - \Phi_0)|_{\bar{h}} \leq C_l \theta^{-\varepsilon}
\]
on \( \bar{M} \times (0, t_0) \) for some \( 0 < t_0 < t_+ \). Here \( \bar{U} := \theta^{-1} \partial_t \). Moreover, \( \bar{D} \) denotes the Levi-Civita connection associated with \( \bar{h} \) and \( q \) denotes the logarithmic volume density, defined by the condition \( e^\rho \mu_h = \mu_g \), where \( \mu_h \) is the volume form on \( M \) induced by \( h \), \( \mu_g \) is the volume form on \( M \) induced by the metric \( \Psi \bar{g} \) on the leaves \( \bar{M} \). If, in addition to the above, \( (M, g) \) is globally hyperbolic and the hypersurfaces \( \Psi(M_t) \) are Cauchy hypersurfaces in \( (M, g) \) for \( t \in (0, t_+) \), then \( (M, g) \) is called a locally crushing globally hyperbolic development with vanishing shift vector field of the initial data.

1.5. The stable quiescent regime. There are results guaranteeing the existence of developments corresponding to non-degenerate quiescent initial data on the singularity for the Einstein-scalar field equations. In the vacuum setting, one example of this is given by [29] Theorem 1.1, p. 2], mentioned in connection with Definition [1]. Moreover, in case condition 4 of Definition [10] is void, the existence of developments is guaranteed, in the real analytic setting, by [6, 22]. In order to see this, it is of interest to relate Definition [10] with the perspective of [6, 22]. The main theorems of [6, 22] (i.e., [6] Theorems 1 and 2, pp. 484–485] and [22] Theorems 10.1 and 10.2, p. 1104) schematically state the following: given real analytic solutions to a so-called velocity dominated version of the equations (or to the Kasner-like Einstein-matter equations), there is a unique real analytic solution to the actual equations whose asymptotics are determined by the velocity dominated solution. In other words, in [6, 22], the velocity dominated solutions should be thought of as initial data on the big bang singularity. In the case of the Einstein-scalar field equations, solutions to the the velocity dominated system are given by \( (0^g, 0^k, 0^\phi) \), where \( 0^g \) is a family of real analytic Riemannian metrics on a manifold \( \bar{M} \), \( 0^k \) is a family of real analytic symmetric covariant 2-tensor fields on \( \bar{M} \) and \( 0^\phi \) is a family of real analytic functions on \( \bar{M} \). The families are parametrised by \( t \in (0, \infty) \) and depend real analytically on \( t \). Needless to say, one could discuss velocity dominated solutions in other regularity classes. The equations that should be satisfied are listed in [6] Subsection 2.2, pp. 482–483]; see [6] (10)–(13), p. 438]. One advantage of the perspective taken in [6, 22] is that the condition of non-degeneracy of \( \mathcal{K} \), cf. Definition [10] does not appear; cf. also Remark [26] below for a further discussion of the issue of non-degeneracy. The following two propositions relate the two perspectives on initial data on the singularity.

**Proposition 18.** Let \( \mathcal{M} = (\bar{M}, \bar{h}, \mathcal{K}, \Phi_0, \Phi_b) \) be non-degenerate quiescent initial data on the singularity for the Einstein-scalar field equations with a vanishing cosmological constant. Denote the eigenvalues of \( \mathcal{K} \) by \( p_A, A = 1, \ldots, n \), and order them so that \( p_1 < \cdots < p_n \). Let \( \mathcal{F}_A, A = 1, \ldots, n \), be an eigenvector field of \( \mathcal{K} \) corresponding to \( p_A \), normalised so that \( |\mathcal{F}_A|_{\bar{h}} = 1 \), and let \( \{ \mathcal{Y}^A \} \) be the basis dual to \( \{ \mathcal{F}_A \} \). Then

\[
(19a) \quad 0^g := \sum A 2^p_A \mathcal{Y}^A \otimes \mathcal{Y}^A,
\]
\[
(19b) \quad 0^k := - \sum A p_A 2^{p_A - 1} \mathcal{Y}^A \otimes \mathcal{Y}^A,
\]
\[
(19c) \quad 0^\phi := \Phi_0 \ln t + \Phi_b
\]
is a velocity dominated solution in the sense of [6] Subsection 2.2, pp. 482–483].

**Remark 19.** The conventions of [6] differ from ours in that a factor of \( 8\pi \) is included in Einstein’s equations in [6]. This means that all the terms involving the scalar field appearing on the right hand sides of [6] (10a), (10b) and (11b), p. 483] should be divided by \( 8\pi \) in order to make the formulae consistent with our conventions.
Remark 20. Note that while the frame $\{\mathcal{A}_a\}$ and co-frame $\{\mathcal{A}^A\}$ may only be local, the $\mathcal{V}^A$ are well defined up to a sign, so that the right hand sides of (19a) and (19b) are globally well defined.

**Proof.** The proof is to be found in Section 4 below. □

In the following proposition, we go in the opposite direction.

**Proposition 21.** Let $(0g, 0k, 0\phi)$ be a velocity dominated solution in the sense of [6], Subsection 2.2, pp. 482–483, where the underlying manifold is $\tilde{M}$ with $\dim \tilde{M} \geq 3$. Adjusting the time interval appropriately, it can then be assumed that $\text{tr}^0 k = -t^{-1}$. Moreover, if $0K$ is the $(1,1)$-tensor field with components $0K^a_h = 0k^a_{bh}$, where the indices are raised with $0g$, then $-t \cdot 0K$ is independent of $t$. Define $\mathcal{X} := -t \cdot 0K$. Then the eigenvalues of $\mathcal{X}$ are real. Assume them to be distinct, label them $p_a, A = 1, \ldots, n$, and order them so that $p_1 < \cdots < p_n$. Let $\mathcal{F}_A, A = 1, \ldots, n$, be a local eigenvector field of $\mathcal{X}$ corresponding to $p_a$, and let $\{\mathcal{V}^A\}$ be the basis dual to $\{\mathcal{F}_A\}$. Then there are functions $\alpha_a, A = 1, \ldots, n$, defined on the subset of $\tilde{M}$ on which $\{\mathcal{F}_A\}$ is defined, which are strictly positive and of the same regularity as $0g$ and $0k$, such that

$$0g = \sum A \alpha^2_A \mathcal{V}^A \otimes \mathcal{V}^A.$$  

Define the Riemannian metric $\mathcal{h}$ on $\tilde{M}$ by

$$\mathcal{h} := \sum A \alpha^2_A \mathcal{V}^A \otimes \mathcal{V}^A;$$

the right hand side is independent of the choice of local frame $\{\mathcal{F}_A\}$ and therefore the left hand side is well defined on $\tilde{M}$. Finally, $\partial_t (p^0 \partial_a \phi) = 0$, so that there are functions $\Phi_a$ and $\Phi_b$, of the same regularity as $0\phi$ and $\partial_0 \phi$, such that $0\phi = \Phi_a \ln t + \Phi_b$. Then $(\tilde{M}, \mathcal{h}, \mathcal{X}, \Phi_a, \Phi_b)$ are non-degenerate quiescent initial data on the singularity for the Einstein-scalar field equations with a vanishing cosmological constant, assuming that $1 + p_1 - p_{n-1} - p_n > 0$.

**Proof.** The proof is to be found in Section 4 below. □

Finally, we appeal to [6, 22] in order to deduce the existence of developments in case condition 4 of Definition 10 is void.

**Theorem 22.** Let $(\tilde{M}, \mathcal{h}, \mathcal{X}, \Phi_a, \Phi_b)$ be non-degenerate quiescent initial data on the singularity for the Einstein-scalar field equations with a vanishing cosmological constant. Assume the manifold and the data to be real analytic. Assume, moreover, that if the eigenvalues of $\mathcal{X}$, say $p_a, A = 1, \ldots, n$, are ordered so that $p_1 < \cdots < p_n$, then $1 + p_1 - p_{n-1} - p_n > 0$. Then there is an open neighbourhood $V$ of $\tilde{M} \times \{0\}$ in $\tilde{M} \times \mathbb{R}$ and a real analytic solution $(M, g, \phi)$ to the Einstein-scalar field equations, where $M := V \cap [\tilde{M} \times (0, \infty)]$, such that the following holds. On $M$, the metric $g$ is given by the right hand side of (13). In this expression, $\mathcal{F}_A, A = 1, \ldots, n$, is an eigenvector field of $\mathcal{X}$ corresponding to $p_a$, normalised so that $|\mathcal{X}_A|_{\mathcal{h}} = 1$, and $\{\mathcal{Y}^A\}$ is the basis dual to $\{\mathcal{F}_A\}$. Moreover, if $K$ is a compact subset of $\tilde{M}$, there are constants $\varepsilon > 0, t_0 > 0$ and $C > 0$ such that

$$\|(K(\cdot, t) - \mathcal{X})\|_{C(K)} \leq Ct^\varepsilon,$$

$$\|(\mathcal{h}(\cdot, t) - \mathcal{h})\|_{C(K)} \leq Ct^\varepsilon,$$

$$\|((\mathcal{U}\phi)(\cdot, t) - \Phi_a)\|_{C(K)} \leq Ct^\varepsilon,$$

$$\|((\phi - \Phi_a g)(\cdot, t) - \Phi_b)\|_{C(K)} \leq Ct^\varepsilon$$

for all $t \leq t_0$. Here $K$, $\mathcal{h}$, $\mathcal{U}$ and $g$ are defined as in Definition 16.

**Remark 23.** Since the arguments are in the real analytic setting, it should be possible to upgrade (21) to $C^t$-estimates. However, we do not derive such estimates here.

**Remark 24.** The results in [6, 22] also include uniqueness statements. Under additional conditions, the constructed solutions are thus unique. We omit the details.
Remark 25. In the course of the proof, it is demonstrated that $t \theta$ converges to 1 at a rate. The metric can thus be also be represented as in (17) with $N = 1$. Moreover, on compact subsets, estimates of the form (18) hold with $k = 0$, though higher order estimates should also hold; cf. Remark 23.

Proof. The proof is to be found in Section 4 below. \qed

Finally, let us make a remark about the degenerate case.

Remark 26. In Definitions 1, 2, 7, 10, 16 and 17, the condition of non-degeneracy is important. This is due to the fact that $\mathcal{O}_{BC}$ is defined in terms of the eigenvector fields corresponding to the distinct eigenvalues and to the fact that the correspondence between $\bar{h}$ and the spacetime metric is based on the fact that the eigenvalues are distinct. However, the notion of a velocity dominated solution, in the sense of [6] Subsection 2.2, pp. 482–483, offers a different perspective in the case of Gaussian foliations. In order to see this, let us focus on the $n = 3$ dimensional case. Ordering the eigenvalues so that $p_1 \leq p_2 \leq p_3$, it is of interest to separate the cases $p_1 = p_2$ and $p_2 = p_3$. If $p_1 > 0$, the condition on the $\mathcal{O}_{BC}^A$ is void. Let us therefore assume that $p_1 \leq 0$. If $p_1 = p_2 \leq 0$, then, due to the conditions on the $p_i$, $p_3 = 1$ and $p_1 = p_2 = 0$. This case is consistent with the flat Kasner solutions (which have Cauchy horizons). It may be of interest to allow this possibility, but it can be expected to be quite different from the situation we are interested in here. We therefore exclude this case. What remains is $p_1 \leq 0$ and $p_2 = p_3$. Then $p_2 = p_3 > 0$. In particular the (2-dimensional) eigenspace corresponding to $p_2$ and $p_3$ is given by the orthogonal complement of the eigenspace of $p_1$. In this setting, one generalisation of the condition $\mathcal{O}_{BC}^1 = 0$ in a neighbourhood of a point $\bar{x}$ such that $p_1(\bar{x}) \leq 0$ is the requirement that there is a neighbourhood, say $V$, of $\bar{x}$ such that the distribution, say $D$, of $\bar{M}$ of rank 2 given by the tangent vectors orthogonal to the eigenspace of $p_1$, is involutive, so that $D$ is integrable (and, in fact, defines a foliation on $V$).

The above considerations lead to the following generalisation in the case of $n = 3$. Fix an asymptotically velocity dominated solution $(\theta^0, \theta^0k^0, \theta^0\phi)$, assume that $\text{tr}^0k = -t^{-1}$ and fix $\mathcal{K} := -t \cdot \theta^0 \bar{K}$; cf. the statement of Proposition 21. Assume the eigenvalues of $\mathcal{K}$ (which are automatically $\leq 1$) to be strictly less than 1. Assume, moreover, that if there is an $\bar{x} \in \bar{M}$ such that the smallest eigenvalue of $\mathcal{K}$ at $\bar{x}$ is $\leq 0$, then there is an open neighbourhood $V$ of $\bar{x}$ such that the distribution, say $D$, of $\bar{M}$ of rank 2 given by the tangent vectors orthogonal to the eigenspace of $p_1$, is involutive. Under these conditions, we expect the existence of a solution to the Einstein-scalar field equations with asymptotics determined by the asymptotically velocity dominated solution as in [6] Theorem 1, p. 484).

1.6. Improving the asymptotics. The correspondence between the developments in Definitions 7, 16 and 17 and the initial data is given by (15), (10) and (18) respectively. However, as pointed out in connection with Definition 7, more detailed information might be needed to deduce uniqueness; cf., e.g., [29] Theorem 1.7, p. 4. In particular, it may be necessary to control the asymptotics of $\theta, \partial_i k^j, \text{etc.}$ It may also be necessary to have more detailed information concerning the rate of convergence. It is therefore of interest to assume that we have developments as in Definitions 7, 16 and 17 and to deduce more detailed information concerning the asymptotics. We do so here. In addition, it turns out to be possible to derive some of the assumptions included in the definitions of a development and of initial data. The statement of the following proposition therefore differ somewhat from the definitions of initial data and of a development.

Theorem 27. Let $3 \leq n \in \mathbb{N}$, $\Lambda \in \mathbb{R}$, $(\bar{M}, \bar{h})$ be a smooth and closed $n$-dimensional Riemannian manifold and $\mathcal{K}$ be a smooth $(1, 1)$-tensor field on $\bar{M}$. Assume that

1. $\mathcal{K}$ is symmetric with respect to $\bar{h}$.
2. $\text{tr} \mathcal{K}^2 \leq 1$.
3. The eigenvalues $p_A$ of $\mathcal{K}$ are distinct.
4. $\mathcal{O}_{BC}^A$ vanishes in a neighbourhood of $\bar{x} \in \bar{M}$ if $1 + p_A(\bar{x}) - p_B(\bar{x}) - p_C(\bar{x}) \leq 0$, where $\mathcal{O}_{BC}^A$ is given by Remark 14.
Assume that there is a smooth solution \((M_0, g, \phi)\) to the Einstein-scalar field equations with a cosmological constant \(\Lambda\) such that \(M_0 = \tilde{M} \times (0, t_+)\) for some \(t_+ > 0\); the mean curvature \(\theta\) of the leaves \(M_t := \tilde{M} \times \{t\}\) is strictly positive; and \(g\) takes the form
\[
g = -N^2 dt \otimes dt + \sum_{A,B} b_{AB} \theta^{-2p_{\max}(A,B)} \mathcal{Y}^A \otimes \mathcal{Y}^B,
\]
where \(\{\mathcal{Y}^A\}\) is the basis dual to \(\{\mathcal{X}_A\}\) introduced in Remark 11; the lapse function \(N\) is strictly positive and the \(p_A\) are ordered so that \(p_1 < \cdots < p_n\). Moreover, \(\theta\) is assumed to diverge uniformly to \(\infty\) as \(t \to 0^+\). Let \(\hat{h}\) be defined by (14). \(K\) be the expansion normalised Weingarten map of the leaves of the foliation and \(\hat{U}\) be the future directed unit normal with respect to \(\hat{g} := \theta^2 g\). Assume (18a) and (18b) to hold. Finally, assume that there is a constant \(C_{\text{rel}}\) such that
\[
|\hat{D}\ln \hat{N}|_{\hat{h}} \leq C_{\text{rel}}
onumber
\]
on \(\hat{M}_0 := \tilde{M} \times (0, t_0]\) for some \(t_0 \in (0, t_+)\), where \(\hat{N} := \Theta N\) and \(\hat{D}\) denotes the Levi-Civita connection induced by \(\hat{h}\), and that there are, for each \(1 \leq k \in \mathbb{N}\), constants \(a_k\) and \(C_k\) such that
\[
|\hat{D}^k \ln \hat{N}|_{\hat{h}} + |\hat{D}^k \phi|_{\hat{h}} + |\hat{D}^k - 1 \hat{U} \ln N|_{\hat{h}} \leq C_k (\ln \theta)^{a_k}
\]
on \(\hat{M}_0\). Let the logarithmic volume density \(\phi\) be defined by the condition \(e^\phi \mu_{\hat{h}} = \mu_{\hat{g}}\), where \(\mu_{\hat{h}}\) is the volume form on \(\tilde{M}\) induced by \(\hat{h}\) and \(\mu_{\hat{g}}\) is the volume form on \(\tilde{M}\) given by the metric \(\hat{g}\) induced on the leaves \(\hat{M}_t\). Let, moreover, \(\{X_A\}\) be the eigenvector field of \(\mathcal{K}\), normalised so that the \(X_A\) are unit vector fields with respect to \(\hat{h}\) and so that \(X_A \to \mathcal{X}_A\). Then \(\tilde{\mu}_{\hat{h}}\) is defined by the condition that \(|X_A|_{\hat{g}} = e^{\bar{\mu}_A}\). With this notation, there is an \(\eta > 0\) and, for each \(k \in \mathbb{N}\), a constant \(C_k\) such that
\[
|\hat{D}^k (\ln \theta + \phi)|_{\hat{h}} + \sum_A |\hat{D}^k (\bar{\mu}_A + p_A \ln \theta)|_{\hat{h}} \leq C_k \theta^{-2\eta}
\]
on \(\hat{M}_0\). Next, there is an \(\eta > 0\) and, for each \(k \in \mathbb{N}\), a constant \(C_k\) such that
\[
|\hat{D}^k (\hat{\mathcal{L}}_U \mathcal{K})|_{\hat{h}} \leq C_k \theta^{-2\eta},
\]
\[
|\hat{D}^k [\mathcal{K} (\mathcal{Y}^A, \mathcal{Y}^B) - p_B \delta^A_B]|_{\hat{h}} \leq C_k \theta^{-2\eta} \min \{1, \theta^{-2(p_B - p_A)}\},
\]
\[
|\hat{D}^k [\hat{\mathcal{L}}_U \mathcal{K}] (\mathcal{Y}^A, \mathcal{Y}^B)|_{\hat{h}} \leq C_k \theta^{-2\eta} \min \{1, \theta^{-2(p_B - p_A)}\}
\]
on \(\hat{M}_0\) for all \(A, B\) (no summation on \(B\)). In addition, if the deceleration parameter \(q\) is defined by the condition that \(\hat{U} (n \ln \theta) = -1 - q\), then there is an \(\eta > 0\) and, for each \(k \in \mathbb{N}\), a constant \(C_k\) such that
\[
|\hat{D}^k [q - (n - 1)]|_{\hat{h}} \leq C_k \theta^{-2\eta}
\]
on \(\hat{M}_0\). Moreover, there are \(\Phi_A, \Phi_b \in C^\infty(\tilde{M})\), an \(\eta > 0\) and, for each \(k \in \mathbb{N}\), a constant \(C_k\) such that
\[
|\hat{D}^k (\hat{U} \phi - \Phi_a)|_{\hat{h}} \leq C_k \theta^{-2\eta},
\]
\[
|\hat{D}^k (\phi - \Phi_a \theta - \Phi_b)|_{\hat{h}} \leq C_k \theta^{-2\eta}
\]
on \(\hat{M}_0\). Finally, \(\text{tr} \mathcal{X} = 1\), \(\text{tr} \mathcal{Z}^2 + \Phi_a^2 = 1\) and \(\text{div}_h \mathcal{X} = \Phi_a d\Phi_b\).

Remark 28. The notation \(\hat{\mathcal{L}}_U \mathcal{K}\) is introduced in [65] Section A.2, pp. 203–204]. However, in the case of a vanishing shift vector field, it can be defined as follows. Let \(\{E_i\}\) be a local frame on \(\tilde{M}\) (in particular, it is independent of \(t\)) and \(\{\omega^i\}\) be the dual frame. Then \(\hat{\mathcal{L}}_U \mathcal{K}^i_j = \hat{U} (\mathcal{K}^i_j)\), where the sub- and superscripts refer to the frame \(\{E_i\}\) and co-frame \(\{\omega^i\}\); cf. [65] (A.3), p. 204].

Remark 29. The frames \(\{\mathcal{X}_A\}\) and \(\{X_A\}\) are only defined locally on \(\tilde{M}\). However, the \(\mathcal{X}_A\) and \(X_A\) are well defined up to a sign (since \(\mathcal{K}\) converges to \(\mathcal{X}\) and \(\mathcal{X}\) has distinct eigenvalues, we can assume \(\mathcal{K}\) to have distinct eigenvalues by restricting \(t_0\), if necessary). This means that \(\tilde{\mu}_A\) is globally well defined. Moreover, the estimates (26b) and (26c) make sense globally.

Remark 30. Note that we do not assume the asymptotic versions of the Hamiltonian and momentum constraints to hold. Moreover, we do not assume (18c) or (18d) to hold. However, we do require that \(\text{tr} \mathcal{Z}^2 \leq 1\). This requirement, when combined with the fact that \(n \geq 3\) and the non-degeneracy, implies that \(p_A < 1\) for all \(A\).
Remark 31. In the statement of the theorem we assume \( \bar{M} \) to be compact. However, by a more technical argument, we expect it should be possible to obtain similar conclusions in the non-compact setting.

Remark 32. The requirement that (24) hold is perhaps the least satisfying assumption. However, it is hard to avoid making assumptions on the lapse function. Here, we bound the first and the last term on the left hand side of (24). What is perhaps more questionable is the bound on the middle term on the left hand side. However, in deriving estimates for the scalar field, we need control of the geometry, in particular on the deceleration parameter. However, such control is only obtained by combining the assumptions with Einstein’s equations. In that context, we need some control of the scalar field. The control necessary could be obtained in different ways. One could also, for example, assume that \(|D^k \hat{U} \phi|_{\bar{h}} \) decays or that \(|D^k \hat{U} \phi|_{\bar{h}} \) does not grow faster than the right hand side of (24). However, that would imply the assumed bound on \( \phi \).

Remark 33. For \( \xi \in \mathbb{R}^n \) and \( \langle \xi \rangle := (1 + |\xi|^2)^{1/2} \).

Remark 34. If \( n = 3 \) and \( \theta_{2A}^3 = 0 \) on \( \bar{M} \), then, due to Lemma [55] below, there are, for each \( p \in \bar{M} \), local coordinates \((V, \tilde{x})\) and strictly positive functions \( \alpha_A, A = 1, 2, 3 \), such that \( p \in V \) and \( \mathcal{J}_A = \alpha_A \mathcal{J}_A \) (no summation), where the \( \mathcal{J}_A \) are given by (58) below. As a consequence, \( \Psi_A = \alpha_A^{-1} \mathcal{J}_A \) (no summation), where the \( \mathcal{J}_A \) are given by (59) below. Moreover,

\[
dx^i = \mathcal{J}^i_A \mathcal{J}_A, \quad \partial_i = \mathcal{J}^A_A \mathcal{J}_A,
\]

where \( \mathcal{J}^2_A, \mathcal{J}_1^2 \) and \( \mathcal{J}^3_A \) are introduced in (58) below. The remaining components of \( \mathcal{J}^2_A \) are determined by \( \mathcal{J}^i_A = 0 \) if \( i < A \) and \( \mathcal{J}^i_A = 1 \) if \( A = i \). Moreover,

\[
\mathcal{J}^2_A = -\mathcal{J}^3_A, \quad \mathcal{J}_1^2 = -\mathcal{J}_1^3, \quad \mathcal{J}^3_A = \mathcal{J}^2_1 \mathcal{J}^3_2 - \mathcal{J}^3_1,
\]

\( \mathcal{J}^i_A = 0 \) if \( A < i \) and \( \mathcal{J}^A_A = 1 \) if \( A = i \). Note also that

\[
dx^i = \mathcal{J}^i_A \Psi_A, \quad \partial_i = \mathcal{J}^A_A \Psi_A,
\]

where \( \mathcal{J}^i_A = \alpha_A^{-1} \mathcal{J}^i_A \) (no summation) and \( \Psi_A = \alpha_A \mathcal{J}^A_A \) (no summation). Combining these observations with (26b) and (26c), it can be verified that

\[
|D^k[K(dx^i, \partial_j) - \mathcal{K}(dx^i, \partial_j)]|_{\bar{h}} \leq C_k \theta^{-2q} \min\{1, \theta^{-2(p_i - p_j)}\},
\]

\[
|\hat{D}^k[(\mathcal{L}_U \mathcal{K})(dx^i, \partial_j)]|_{\bar{h}} \leq C_k \theta^{-2q} \min\{1, \theta^{-2(p_i - p_j)}\}
\]
on \( M_a \) for all \( i, j \).

Proof. The proof is to be found in Subsection 5.5 below.

In the case of a Gaussian foliation, there is a similar statement. However, the assumptions are slightly different.

Theorem 35. Let \( 3 \leq n \in \mathbb{N}, \Lambda \in \mathbb{R}, (\bar{M}, \bar{h}) \) be a smooth and closed \( n \)-dimensional Riemannian manifold, \( \mathcal{K} \) be a smooth \((1,1)\)-tensor field and \( \Phi_a \) a smooth function on \( \bar{M} \). Assume that

(1) \( \mathcal{K} \) is symmetric with respect to \( \bar{h} \).
(2) \( \tr \mathcal{K}^2 + \Phi_a^2 = 1 \).
(3) The eigenvalues \( p_A \) of \( \mathcal{K} \) are distinct.
(4) \( \mathcal{K} \) vanishes in a neighbourhood of \( \bar{x} \in \bar{M} \) if \( 1 + p_A(\bar{x}) - p_B(\bar{x}) - p_C(\bar{x}) \leq 0 \), where \( \mathcal{K}^A_{BC} \) is given by Remark [14].

Assume that there is a smooth solution \((M_a, g, \phi)\) to the Einstein-scalar field equations with a cosmological constant \( \Lambda \) such that \( M_a = [0, t_+] \times (0, t_+) \times \mathbb{R}^n \) for some \( t_+ > 0 \); the mean curvature \( \theta \) of the leaves \( M_t := \bar{M} \times \{t\} \) satisfies \( \theta \geq \theta_0 \) on \( M_a \) for some \( 0 < \theta_0 < 1 \), \( \theta_0 \) is additionally required to satisfy \( \theta_0 > [2\Lambda/(n-1)]^{1/2} \) in case \( \Lambda > 0 \). Assume, in addition, that \( g \) takes the form

\[
g = -dt \otimes dt + \sum_{A,B} B_{AB} t^{2p_{\max}(A, B)} \Psi_A \otimes \Psi_B,
\]

where \( \{\Psi_A\} \) is the basis dual to \( \{\mathcal{J}^A_A\} \) introduced in Remark [14] and the \( p_A \) are ordered so that \( p_1 < \cdots < p_n \). Let \( \bar{h} \) be defined by [14], \( K \) be the expansion normalised Weingarten map of the
leaves of the foliation and \( \hat{U} \) be the future directed unit normal with respect to \( \hat{g} := \theta^2 g \). Assume that (16a)–(16c) hold with \( \varepsilon \) normalised so that the X given by the metric \( \hat{h} \) leaves of the foliation and is defined by the condition that \( \bar{K}_a \) for (29a)

\[
\| \ln \theta + \phi \|_{C^0(\hat{M})} + \sum A ||\tilde{\mu}_A(\cdot, t) - p_A \ln t ||_{C^0(\hat{M})} \leq C_k t^{2\eta},
\]

for \( t \leq t_0 \). Next, there is an \( \eta > 0 \) and, for each \( k \in \mathbb{N} \), a constant \( C_k \) such that

\[
(30a) \quad \| (\hat{L}_U \mathcal{K})(\cdot, t) \|_{C^0(\hat{M})} \leq C_k t^{2\eta},
\]

\[
(30b) \quad \| \mathcal{K}(\mathcal{A}, \mathcal{B}) - p_B \delta^A_B \|_{C^0(\hat{M})} \leq C_k t^{2\eta} \min \{ 1, t^{2(p_B - p_A)} \},
\]

\[
(30c) \quad \| (\hat{L}_U \mathcal{K})(\mathcal{A}, \mathcal{B})(\cdot, t) \|_{C^0(\hat{M})} \leq C_k t^{2\eta} \min \{ 1, t^{2(p_B - p_A)} \},
\]

for \( t \leq t_0 \) and all \( A, B \) (no summation on \( B \)). In addition, if the deceleration parameter \( q \) is defined by the condition that \( \hat{U}(n \ln \theta) = -1 - q \), then there is an \( \eta > 0 \) and, for each \( k \in \mathbb{N} \), a constant \( C_k \) such that

\[
\| q(\cdot, t) - (n - 1) \|_{C^0(\hat{M})} \leq C_k t^{2\eta}
\]

for \( t \leq t_0 \). Moreover, there is a \( \Phi_0 \in C^\infty(\hat{M}) \), an \( \eta > 0 \) and, for each \( k \in \mathbb{N} \), a constant \( C_k \) such that, in addition to (16c),

\[
\| \phi(\cdot, t) - \Phi_0 \phi(\cdot, t) - \Phi_0 \|_{C^0(\hat{M})} \leq C_k t^{2\eta}
\]

for \( t \leq t_0 \). Finally, \( \text{tr} \mathcal{K} = 1 \) and \( \text{div} \mathcal{K} = \Phi_0 \partial \Phi_0 \).

Remark 36. Comments analogous to Remarks 28, 29, and 31 are equally relevant here.

Remark 37. The assumptions of this theorem are stronger than those of Theorem 27 in the sense that we assume \( \text{tr} \mathcal{K}^2 + \phi_2 = 1 \) to be satisfied and \( \hat{U} \phi \) to converge to \( \Phi_0 \). On the other hand, they are weaker in the sense that we assume neither (23) nor (24) to hold. Note also that we obtain the conclusions (29b) concerning the asymptotics of \( \theta \), even though the assumptions involve expansion-normalised quantities.

Remark 38. Assuming, in addition to the conditions of the theorem, that \( n = 3 \) and \( \theta^1_{23} = 0 \) on \( \hat{M} \), then, combining the conclusions of the theorem with the observations of Remark 34 yields

\[
t|\hat{B}^k (\hat{K}_j^i - t^{-1} \mathcal{K}_j^i)|_h \leq C_k t^{2\eta} \min \{ 1, t^{2(p_B - p_A)} \},
\]

\[
t^2|\hat{D}^k \partial_i (\hat{K}_j^i - t^{-1} \mathcal{K}_j^i)|_h \leq C_k t^{2\eta} \min \{ 1, t^{2(p_B - p_A)} \}
\]

with respect to the local coordinates introduced in Remark 34. Here \( \mathcal{K}_j^i \) denotes the components of \( \mathcal{K} \) with respect to the frame \{ \partial_i \} and its dual frame \{ \partial_i^* \}, and similarly for \( \hat{K} \). These estimates should be compared with (6).

\[
\text{Proof.} \quad \text{The proof of this statement is given in Subsection 6.1 below.}
\]

1.7. From convergence to data on the singularity. It is of great interest to characterise the set of vacuum solutions that converge to the Kasner circle. Moreover, for reasons clarified in Subsection 2.2 below, the non-degenerate subset of the Kasner circle is of greatest interest (i.e., the subset on which the \( p_i \) are distinct). More generally, the same question can be asked concerning the Kasner disc. The situation we have in mind is thus that \( \hat{L}_U \mathcal{K} \) (cf. Remark 25) converges to zero; that \( \hat{U}^2 \phi \) converges to zero; that \( \mathcal{K} \) converges to a limit \( \mathcal{K} \); that \( \hat{U} \phi \) converges to a limit \( \Phi_0 \); and that \( \phi^2 + \text{tr} \mathcal{K}^2 = 1 \). However, in order to derive meaningful estimates, we need to choose a norm with respect to which we assume convergence. We should also specify a rate. In the case of a Gaussian foliation, one could use powers of \( t \) to specify the rate. However, this does not work more generally. Another option is negative powers of the mean curvature. However, this quantity
involves one derivative of the metric. On the other hand, in the situations discussed above, the logarithmic volume density \( \varrho \) is comparable in size to the logarithm of the mean curvature; cf., e.g., (25) and (29). This quantity does not involve any derivatives of the metric and can serve as a substitute for the logarithm of the mean curvature. In what follows, we therefore use \( \varrho \) to specify a rate. For this reason, we, at the very minimum, need to require \( \varrho \) to tend to \(-\infty\) uniformly. This leads to the following definition.

**Definition 39.** Let \( 3 \leq n \in \mathbb{N} \), \( M \) be a closed \( n \)-dimensional manifold, \( 0 < t_+ \in \mathbb{R}, I = (0, t_+) \) and \( \bar{M} := M \times \{t\} \). Let \( \Lambda \in \mathbb{R} \). Assume \((M, g, \varphi)\) to be a time oriented Lorentz manifold and \( \phi \) to be a smooth function on \( M \), solving the Einstein-scalar field equations with a cosmological constant \( \Lambda \). Assume the leaves \( \bar{M}_t := \bar{M} \times \{t\} \) of the foliation to be spacelike and denote the future directed unit normal by \( U \). Define the lapse function \( N \) and the shift vector field \( \chi \) by \( \partial_t = NU + \chi \). Assume \( N > 0 \) and \( \chi = 0 \). Assume, moreover, that there is a \( 0 < \theta_0 \in \mathbb{R} \) such that the mean curvature \( \theta \) of the leaves of the foliation satisfies \( \theta \geq \theta_0 \). In case \( \Lambda > 0 \), assume, in addition, that \( \theta_0 > [2\Lambda/(n - 1)]^{1/2} \). Fix a smooth reference metric \( \bar{g}_{\text{ref}} \) on \( \bar{M} \) and define the logarithmic volume density \( \varrho \) by the requirement that \( e^\varrho g_{\bar{g}_{\text{ref}}} = \mu \bar{g} \), where \( \bar{g} \) is the metric induced on the leaves of the foliation. Assume, finally, that \( \varrho \) diverges uniformly to \(-\infty\) as \( t \to 0 \) in the sense that for every \( C \in \mathbb{R} \), there is a \( T \in I \) such that \( \varrho(x, t) \leq C \) for all \( t \leq T \) and \( x \in \bar{M} \). Then \((M, g, \varphi)\) is said to be a solution to the Einstein-scalar field equations with a cosmological constant \( \Lambda \) and a uniform volume singularity at \( t = 0 \).

Turning to the formal requirement of decay of \( \hat{L}_U \mathcal{K} \), we are interested in solutions to the Einstein-scalar field equations with a cosmological constant \( \Lambda \) and a uniform volume singularity at \( t = 0 \) (in the sense of Definition 39) such that there is a function \( \epsilon : \bar{M} \to (0, 1) \) satisfying \( \epsilon \geq \epsilon_0 \) for some \( \epsilon_0 > 0 \), a \( k \in \mathbb{N} \) and constants \( a_k \) and \( C_k \) such that

\[
\sum_{l=0}^k |D^l \hat{L}_U \mathcal{K}|_{\bar{g}_{\text{ref}}} \leq C_k(g)^{a_k} e^{2\epsilon \varrho}
\]

on \( M \), where \( \hat{D} \) is the Levi-Civita connection of \( \bar{g}_{\text{ref}} \). In what follows, \( a_k \) and \( C_k \) will change from line to line, but the function \( \epsilon \) remains fixed. In analogy with (31), we demand that

\[
\sum_{l=0}^k |\hat{D}^l \hat{U}^2 \phi|_{\bar{g}_{\text{ref}}} \leq C_k(g)^{a_k} e^{2\epsilon \varrho}
\]

on \( M \).

In case (31) holds with \( k = 0 \), it can be deduced that \( \mathcal{K} \) converges uniformly to a continuous limit \( \mathcal{K} \); cf. Lemma 88 below. As mentioned previously, we are here interested in the non-degenerate setting. We therefore assume the eigenvalues of \( \mathcal{K} \) to be distinct. Due to Lemma 88 below, it is no restriction to, in this situation, make the following assumption.

**Definition 40.** Let \((M, g, \varphi)\) be a solution to the Einstein-scalar field equations with a cosmological constant \( \Lambda \) and a uniform volume singularity at \( t = 0 \). Let \( \mathcal{K} \) denote the associated expansion normalised Weingarten map and denote the eigenvalues of \( \mathcal{K} \) by \( \ell_A, A = 1, \ldots, n \). Assume that (31) holds with \( k = 0 \). Assume, moreover, that there is an \( \epsilon_{\text{nd}}>0 \) such that \( \min\{\ell_A - \ell_B\} \geq \epsilon_{\text{nd}} \) on \( M \), that \( \ell_1 < \cdots < \ell_n \) and that there are global eigenvector fields \( X_A \) such that \( \mathcal{K} X_A = \ell_A X_A \) (no summation) and \( |X_A|_{\bar{g}_{\text{ref}}} = 1 \). Then the **standard assumptions** are said to be satisfied.

**Remark 41.** Under the standard assumptions, the expansion normalised Weingarten map \( \mathcal{K} \) converges uniformly to a limit \( \mathcal{K} \). Moreover, \( \mathcal{K} \) has distinct eigenvalues, denoted \( p_1 < \cdots < p_n \), and there is a frame of eigenvector fields \( \{Y^A\} \) of \( \mathcal{K} \) (with \( \mathcal{K} Y^A = p_A Y^A \) (no summation)), with dual frame \( \{Y^A\} \), such that \( X_A \) and \( Y^A \) converge uniformly to \( X_A \) and \( Y^A \) respectively, where \( \{Y^A\} \) is the dual basis of \( \{X_A\} \). In fact, (129) and (133) below hold.

Below we use the following terminology.

**Definition 42.** Given that the standard assumptions hold, let \( \{X_A\} \), \( \{Y^A\} \), \( \{X_A\} \) and \( \{Y^A\} \) be given by Definition 40 and Remark 41. Define \( X^B_A \) and \( Y^B_A \) by the relations \( X_A = X^B_A X_B \) and \( Y^A = Y^B_A Y^B \). Finally, define \( \mu_A \) and \( \mu_A \) by \( |X_A|_{\bar{g}} = e^{\mu_A} \) and \( \mu_A := \mu_A + \ln \theta \).

Given these assumptions, a \( C^0 \)-version of the form (17) can be derived with \( \theta \) replaced by \( e^{-\epsilon} \).
Theorem 43. Assume the standard assumptions to be satisfied and let $p_A$, $\{\mathcal{Y}_A\}$, $\{\mathcal{Y}^A\}$, $X^B_A$, $Y^B$ and $\mu_A$ be given by Remark 41 and Definition 42. Then

$$g = -N^2 dt \otimes dt + \sum_{A,B} b_{AB}^{2p_{\max}(A,B)} e^{\mathcal{Y}^A \otimes \mathcal{Y}^B}$$

on $M$. Moreover, there are constants $a$ and $C$, $a$, $t_0 \in (0, t_+)$ and continuous functions $r_A$ on $\bar{M}$ such that

$$|b_{AB} - e^{2r_A \delta_{AB}}| \leq C(\theta)^a e^{2e^\theta}$$

(no summation on $A$) on $M_0 := \bar{M} \times (0, t_0]$. Finally,

$$|\mu_A - p_A e^{-r_A}| \leq C(\theta)^a e^{2e^\theta},$$

$$|X^A_B - \delta^A_B| + |Y^A_B - \delta^A_B| \leq C(\theta)^a e^{2e^\theta} \min\{1, e^{2(p_B - p_A)\theta}\},$$

$$|\mathcal{K}(\mathcal{Y}^A, \mathcal{Y}^B) - p_B \delta^A_B| \leq C(\theta)^a e^{2e^\theta} \min\{1, e^{2(p_B - p_A)\theta}\},$$

$$|\mathcal{L}_{\bar{U}} \mathcal{K}(\mathcal{Y}^A, \mathcal{Y}^B)| \leq C(\theta)^a e^{2e^\theta} \min\{1, e^{2(p_B - p_A)\theta}\}$$

(no summation on $B$) on $M_0$.

Remark 44. Assume, in addition, that (32) holds for $k = 0$, so that $\bar{U} \phi$ converges to a limit, say $\Phi_0$; cf. Lemma 84. Assuming, moreover, that $\Phi_0^A + \text{tr}\mathcal{K}^2 = 1$, it can be deduced that there are constants $C$ and $a$ and a continuous function $r_A$ such that

$$|\ln \theta + \phi - \theta| \leq C(\theta)^a e^{2e^\theta}$$

on $M$; cf. Lemma 90 below. In particular, we can thus, effectively, replace $\theta$ with $-\ln \theta$ in the conclusions of the theorem.

Proof. The proof is to be found at the end of Subsection 7.1 prior to Lemma 90.

In order to obtain estimates for higher order derivatives, we also need to impose conditions on the lapse function. In practice, it turns out to be convenient to impose conditions on $\bar{N} = \theta \bar{N}$.

Definition 45. If the standard assumptions hold, cf. Definition 40, (31) holds for some $1 \leq k \in \mathbb{N}$; and there are constants $C_k$ and $a_k$ such that

$$\sum_{l=1}^k |\bar{D}^l \ln \bar{N}|_{\bar{g}_{\text{ref}}} \leq C_k(\theta)^{a_k}$$

on $M$, then the $k$-standard assumptions are said to hold.

Theorem 46. Let $1 \leq k \in \mathbb{N}$ and assume the $k$-standard assumptions to hold; cf. Definition 45. Let $\mathcal{X}$, $p_A$, $\{\mathcal{X}_A\}$, $\{\mathcal{X}^A\}$, $X^B_A$, $Y^B$, $\mu_A$ and $\bar{\mu}_A$ be given by Remark 41 and Definition 42. Then $\bar{\mu}_A$ and $\mu_A$ are $C^\infty$ and $p_A$, $\{\mathcal{X}_A\}$, $\{\mathcal{X}^A\}$, $X^B_A$ and $Y^B$ are $C^k$. Moreover,

$$g = -N^2 dt \otimes dt + \sum_{A,B} b_{AB}^{2p_{\max}(A,B)} e^{\mathcal{Y}^A \otimes \mathcal{Y}^B}$$

and there are constants $C_k$ and $a_k$ and a $t_0 \in (0, t_+)$ such that

$$|\bar{D}^l(b_{AB} - e^{2r_A \delta_{AB}})|_{\bar{g}_{\text{ref}}} \leq C_k(\theta)^{a_k} e^{2e^\theta}$$

(no summation on $A$) on $M_0 := \bar{M} \times (0, t_0]$ for all $A, B \in \{1, \ldots, n\}$ and all $0 \leq l \leq k$. Here $r_A$ are $C^k$ functions on $M$ that coincide with the functions $r_A$ obtained in Theorem 43. Moreover,

$$|\bar{D}^l(\bar{\mu}_A - p_A e^{-r_A})|_{\bar{g}_{\text{ref}}} \leq C_k(\theta)^{a_k} e^{2e^\theta},$$

$$|\bar{D}^l(X^A_B - \delta^A_B)|_{\bar{g}_{\text{ref}}} + |\bar{D}^l(Y^A_B - \delta^A_B)|_{\bar{g}_{\text{ref}}} \leq C_k(\theta)^{a_k} e^{2e^\theta} \min\{1, e^{2(p_B - p_A)\theta}\},$$

$$|\bar{D}^l(\mathcal{K}(\mathcal{Y}^A, \mathcal{Y}^B) - p_B \delta^A_B)|_{\bar{g}_{\text{ref}}} \leq C_k(\theta)^{a_k} e^{2e^\theta} \min\{1, e^{2(p_B - p_A)\theta}\},$$

$$|\bar{D}^l(\bar{L}_{\bar{U}} \mathcal{K}(\mathcal{Y}^A, \mathcal{Y}^B))|_{\bar{g}_{\text{ref}}} \leq C_k(\theta)^{a_k} e^{2e^\theta} \min\{1, e^{2(p_B - p_A)\theta}\}$$

(no summation on $B$) on $M_0$ for all $A, B \in \{1, \ldots, n\}$ and all $0 \leq l \leq k$.

Assume, in addition to the above, that (32) holds. Then there are $C^k$-functions $\Phi_a$ and $\Phi_b$ on $\bar{M}$ such that

$$|\bar{D}^l(\bar{U} - \phi)|_{\bar{g}_{\text{ref}}} + |\bar{D}^l(\phi - \Phi_a e - \Phi_b)|_{\bar{g}_{\text{ref}}} \leq C_k(\theta)^{a_k} e^{2e^\theta}$$
on $M_0$ for all $0 \leq l \leq k$. Assume that $\Phi_0$ satisfies $\Phi_0^2 + \text{tr} \kappa^2 = 1$. Then there is a function $r_0 \in C^k(M)$ such that

$$|\bar{D}(\ln \theta + q - r_0)|_{\bar{g}_{\text{ref}}} \leq C_k(\theta)^{\alpha_k} e^{2r_0}$$

on $M_0$ for all $0 \leq l \leq k$. Thus (37) and (38) can be rewritten

$$g = - N^2 dt \otimes dt + \sum_{A,B} b_{AB} \theta^{-2p_{\max}(A,B)} \gamma^A \otimes \gamma^B,$$

(42a)

$$|\bar{D}(\bar{b}_{AB} - e^{2\bar{r}_A} \delta_{AB})|_{\bar{g}_{\text{ref}}} \leq C_k(\ln \theta)^{\alpha_k} \theta^{-2\epsilon},$$

(42b) \hspace{1cm} (no summation on $A$) on $M_0$ for $0 \leq l \leq k$. Here $\bar{r}_A = r_A + p_A r_0$ and

$$\bar{b}_{AB} = b_{AB} \theta^{-2p_{\max}(A,B)}(\theta + \ln \theta),$$

where $b_{AB}$ are the functions appearing in (37). Next,

$$|\bar{D}(\bar{\mu}_A - p_A \theta - r_A)|_{\bar{g}_{\text{ref}}} \leq C_k(\ln \theta)^{\alpha_k} \theta^{-2\epsilon},$$

(44a)

$$|\bar{D}(\theta - r_A - \theta)|_{\bar{g}_{\text{ref}}} \leq C_k(\ln \theta)^{\alpha_k} \theta^{-2\epsilon},$$

(44b)

$$|\bar{D}(X^A - \delta^A_B)|_{\bar{g}_{\text{ref}}} \leq C_k(\ln \theta)^{\alpha_k} \theta^{-2\epsilon},$$

(44c)

$$|\bar{D}(\mathcal{K}(\gamma^A, \bar{\gamma}^A) - py^A \theta^B)|_{\bar{g}_{\text{ref}}} \leq C_k(\ln \theta)^{\alpha_k} \theta^{-2\epsilon},$$

(44d)

$$|\bar{D}([\hat{\mathcal{U}}_A, \bar{\gamma}^A] \bar{\gamma}^A)|_{\bar{g}_{\text{ref}}} \leq C_k(\ln \theta)^{\alpha_k} \theta^{-2\epsilon},$$

(44e)

(45)

on $M_0$ for $0 \leq l \leq k$.

If, in addition, $N = 1$, then

$$|\bar{D}(\partial_t - (n - 1))|_{\bar{g}_{\text{ref}}} \leq C_k(\ln \theta)^{\alpha_k} \theta^{-2\epsilon},$$

(46)

Here

$$|\bar{D}(\partial_t - 1)|_{\bar{g}_{\text{ref}}} \leq C_k(\ln \theta)^{\alpha_k} \theta^{-2\epsilon},$$

(47a)

$$|\bar{D}(\delta_{AB} - e^{2\bar{r}_A} \delta_{AB})|_{\bar{g}_{\text{ref}}} \leq C_k(\ln \theta)^{\alpha_k} \theta^{-2\epsilon},$$

(47b) \hspace{1cm} (no summation on $A$ or $B$) on $M_0$ for $0 \leq l \leq k$. Here

$$c_{AB} := \bar{b}_{AB}(\theta)^{-2p_{\max}(A,B)},$$

where $b_{AB}$ is given by (37).

**Remark 47.** When (41) is satisfied, $\theta$ can be replaced by $- \ln \theta + r_0$ in (40), (44a) and (44b). In other words, modifying $\Phi_0$ etc., $\theta$ can effectively be replaced by $- \ln \theta$.

**Remark 48.** As in the previous remark, when (46) holds, $\bar{\gamma}$ can, effectively, be replaced by $t^{-1}$ in (44).

**Proof.** The proof is to be found at the end of Subsection 7.2. 

One of the main motivations behind the notion of initial data on the singularity introduced in this article is the expectation that it can be used to parametrise the set of solutions to the Einstein-scalar field equations with quiescent singularities, in the non-degenerate setting. In order to justify this expectation, it is necessary to extract initial data on the singularity from the asymptotics of a convergent solution. Moreover, it is necessary to prove the desired relation between the solution and the extracted initial data. In other words, assuming a solution to be convergent, we need to reproduce the conditions appearing in Definitions 10 and 17. Assume, to this end, that the conditions of Theorem 46 are satisfied, except for the requirement that $N = 1$ (in the last part of the theorem). Using the notation of Theorem 46, introduce the $C^k$ Riemannian metric

$$\bar{h} := \sum_A e^{2\bar{r}_A} \gamma^A \otimes \gamma^A.$$
Before proceeding, note that the \( \rho \) appearing in the statement of Theorem 46 differs from the \( \rho \) appearing in Definition 17; the difference is due to the use of different reference metrics in the respective definitions. Here we use \( \rho \) to denote the object appearing in the statement of Theorem 46, but introduce \( \bar{\rho} \) to denote the logarithmic volume density defined using \( \bar{h} \) as a reference metric; i.e., \( e^{\bar{\rho}}\mu_{\bar{h}} = \mu_{\bar{\rho}} \). Next, in order to be consistent with the representation of the metric, introduce \( \bar{\mathcal{F}}_A = e^{-\bar{\rho}}\mathcal{F}_A \) and let \( \{\bar{\mathcal{F}}_A\} \) be the frame dual to \( \{\mathcal{F}_A\} \). Then \( \{\mathcal{F}_A\} \) is an orthonormal frame with respect to \( \bar{h} \) and \( g \) can be written

\[
(49) \quad g = -N^2 dt \otimes dt + \sum_{A,B} \bar{b}_{AB} \theta^{-2p_{\max(A,B)}} \bar{\mathcal{F}}_A \otimes \bar{\mathcal{F}}_B,
\]

where \( \bar{b}_{AB} = e^{-\bar{\rho}}b_{AB} \) (no summation) are \( C^k \) functions satisfying

\[
(50) \quad |D^l(b_{AB} - \delta_{AB})|_{\bar{g}_{\text{ref}}} \leq C_k (\ln \theta)^{a_k} \theta^{-2\epsilon}
\]
on \( M_0 \) for \( 0 \leq l \leq k \); cf. (42b). Combining (49) and (50) with arguments similar to the ones given at the beginning of Subsection 5.1 below yields the conclusion that \( \bar{\rho} + \ln \theta \) converges to zero; cf. (77) below. Since \( \bar{\rho} - \rho \) is constant, and since (41) holds, we conclude that \( \bar{\rho} = \rho - r_B \). Recalling that (40) holds, where the functions \( \Phi_a \) and \( \Phi_b \) are \( C^k \), we conclude that

\[
(51) \quad \sum_{i=0}^k |D^i(\bar{\Phi}(\phi) - \bar{\Phi}_a)|_{\bar{g}_{\text{ref}}} + \sum_{i=0}^k |D^i(\bar{\Phi} - \bar{\Phi}_a \bar{\Phi} - \bar{\Phi}_b)|_{\bar{g}_{\text{ref}}} \leq C_k (\ln \theta)^{a_k} \theta^{-2\epsilon}
\]
on \( M_0 \), where \( \bar{\Phi}_a := \Phi_a \) and \( \bar{\Phi}_b := \Phi_b + r_B \Phi_a \). Next, \( \mathcal{K} \) is \( C^k \) due to Lemma 91 below. With these definitions, the data at the singularity are \((\bar{M}, \bar{h}, \mathcal{K}, \bar{\Phi}_a, \bar{\Phi}_b)\). Finally, as in Definition 17 we define the \( C^k \)-family \( \bar{h} \) of Riemannian metrics on \( M \) by

\[
(52) \quad \bar{h} := \sum_{A,B} \bar{b}_{AB} \bar{\mathcal{F}}_A \otimes \bar{\mathcal{F}}_B.
\]

**Theorem 49.** Assume the 2-standard assumptions, cf. Definition 45, to hold and [32] to hold with \( k = 1 \). Using the notation introduced prior to the statement of the theorem, assume that \( \Phi_a^2 + \text{tr} \mathcal{K}^2 = 1 \). Then, if \( \bar{x} \in \bar{M} \) is such that \( 1 + p_A(\bar{x}) - p_B(\bar{x}) - p_C(\bar{x}) \leq 0 \), there is a neighbourhood \( \mathcal{V} \) of \( \bar{x} \) such that \( \mathcal{F}(\bar{M}, \bar{h}, \mathcal{K}, \bar{\Phi}_a, \bar{\Phi}_b) \) is 0 on \( V \).

Assume, in addition, that the \( k \)-standard assumptions hold with \( k \geq 2 \) and that [32] holds for this \( k \). Define \((\bar{M}, \bar{h}, \mathcal{K}, \bar{\Phi}_a, \bar{\Phi}_b)\) as described prior to the statement of the theorem, let \( \mathcal{K} \) be the expansion normalised Weingarten map and let \( \bar{h} \) be defined by (52). Then \( \text{tr} \mathcal{K} = 1 \); \( \mathcal{K} \) is symmetric with respect to \( \bar{h} \); \( \text{tr} \mathcal{K}^2 + \Phi_a^2 = 1 \); \( \text{div} \mathcal{K} \mathcal{K} = \Phi_a d \Phi_b \); the eigenvalues of \( \mathcal{K} \) are distinct; and \( \mathcal{F}|_{\mathcal{B}_C} = 0 \) vanishes in a neighbourhood of \( \bar{x} \) if \( 1 + p_A(\bar{x}) - p_B(\bar{x}) - p_C(\bar{x}) \leq 0 \). Moreover, there are constants \( C_k \) and \( a_k \) and a \( t_0 \in (0, t_\ast) \) such that (44) and

\[
(53a) \quad \sum_{i=0}^k |D^i(\mathcal{K} - \mathcal{K})|_{\bar{g}_{\text{ref}}} + \sum_{i=0}^k |D^i(\bar{h} - \bar{h})|_{\bar{g}_{\text{ref}}} \leq C_k (\ln \theta)^{a_k} \theta^{-2\epsilon},
\]

\[
(53b) \quad \sum_{i=0}^k |D^i(\bar{\Phi}(\phi) - \bar{\Phi}_a)|_{\bar{g}_{\text{ref}}} + \sum_{i=0}^k |D^i(\bar{\Phi} - \bar{\Phi}_a \bar{\Phi} - \bar{\Phi}_b)|_{\bar{g}_{\text{ref}}} \leq C_k (\ln \theta)^{a_k} \theta^{-2\epsilon},
\]

hold on \( M_0 := \bar{M} \times (0, t_0) \). In particular, since (41) holds, \( \theta \) converges uniformly to \( \infty \) as \( t \to 0^+ \).

**Remark 50.** The theorem demonstrates the necessity of the conditions appearing in Definition 10. Note also that we deduce the form (17) of the metric and the convergence of \( \bar{h} \) to \( h \).

**Proof.** The proof is to be found at the beginning of Subsection 7.3 \( \square \)

Finally, let us return to the original setting: a Gaussian foliation and Einstein’s vacuum equations in 3 + 1-dimensions; cf. Subsection 1.1.

**Theorem 51.** Assume that the \( k \)-standard assumptions hold with \( n = 3 \) and \( a \leq k \in \mathbb{N} \). Assume, moreover, that \( \text{tr} \mathcal{K}^2 = 1 \), that \( N = 1 \) and that the scalar field vanishes. Then \( \mathcal{K} \) converges to its limit \( \mathcal{K} \) in \( C^k \) and the eigenvalues \( p_A \) of \( \mathcal{K} \) satisfy \( p_1 < p_2 < p_3 \), \( \sum_i p_i = 1 \) and \( \sum_i p_i^2 = 1 \). Fix \( \bar{x} \in \bar{M} \). Assuming \( \mathcal{K} \) to be \( C^\infty \), there are local coordinates \((V, \bar{x})\) in a neighbourhood of \( \bar{x} \) such that

\[
g = -dt \otimes dt + \sum_{i,j} d\bar{x}_i t^{2p_{\max(i,j)}} d\bar{x}_i \otimes d\bar{x}_j
\]
on $V \times (0, t_+)$, $\bar{M}$ such that the closure of $W$ is compact and contained in $V$. There are functions $\bar{d}_{ij}$ on $V$, symmetric in $i$ and $j$, constants $C_k$ and $0 < \eta \in \mathbb{R}$ and a $t_0 \in (0, t_+)$ such that
\begin{equation}
\sum_{l \leq k} |\nabla (\bar{d}_{ij} - \bar{d}_{ij})|_{\bar{g}_{\text{ref}}} \leq C_k t^{2\eta}
\end{equation}
on $W \times (0, t_0)$, where $\nabla$ is the Levi-Civita connection associated with $\bar{g}_{\text{ref}}$, $\bar{d}_{ij} > 0$ (no summation),
\begin{equation}
\bar{d}_{12} = \frac{\bar{d}_{22}}{p_2 - p_1} \mathcal{K}^2, \quad \bar{d}_{23} = \frac{\bar{d}_{33}}{p_3 - p_2} \mathcal{K}^3, \quad \bar{d}_{13} = \frac{\bar{d}_{33}}{p_3 - p_1} \left( \mathcal{K}^3 + \frac{1}{p_3 - p_2} \mathcal{K}^2 \mathcal{K}^3 \right).
\end{equation}
Here $\mathcal{K}^i = \mathcal{K}(d^2, \partial_j)$, $\mathcal{K}^i = p_i$ (no summation) and $\mathcal{K}^i = 0$ if $i < j$. Moreover,
\begin{equation}
t \sum_{l \leq k} \left| \nabla (\bar{K}^{i} - t^{-1} \mathcal{K}^i) \right|_{\bar{g}_{\text{ref}}} \leq C_k (\ln t)^{a_k} t^{2\eta} \min \{ 1, t^{2(p_j - p_i)} \},
\end{equation}
\begin{equation}
t^2 \sum_{l \leq k} \left| \nabla \partial_l (\bar{K}^{i} - t^{-1} \mathcal{K}^i) \right|_{\bar{g}_{\text{ref}}} \leq C_k (\ln t)^{a_k} t^{2\eta} \min \{ 1, t^{2(p_j - p_i)} \}
\end{equation}
on $W \times (0, t_0)$. Finally, (3) holds on $V$ with $c_{ij}$ replaced by $\bar{d}_{ij}$ and $\kappa_{ij} = -\mathcal{K}^j$.

**Remark 52.** It is of course not necessary to assume $\mathcal{K}$ to be smooth. The reason we do so here is to avoid technicalities. The interested reader is encouraged to generalise Lemma 55 below to the case that $\mathcal{K}$ is $C^k$. This means that the coordinate system constructed in Lemma 55 below has finite regularity, and one has to trace the consequences of this for the following estimates.

**Remark 53.** The conclusions should be compared with the results of [29], in particular, [29] Theorem 1.1, pp. 2-3 and [29] Theorem 1.7, p. 4; see also Subsection 1.1. Note, in addition, that the assumption (56a) in the current setting corresponds to the requirement that the relative spatial variation of $\theta$ does not grow faster than polynomially in $\varrho$. On the other hand, the conclusion (56a) implies that the relative spatial variation of $\theta$ decays exponentially in $\varrho$.

**Proof.** The proof is to be found at the end of Subsection 7.3.

1.8. **Conclusions and outlook.** Due to the results of this section, there are reasons to conjecture that initial data as in Definition 10 parametrise the stable manifolds of the corresponding subsets of the Kasner disc. The reason for this is that by assuming the solution to converge to a non-degenerate subset of the Kasner disc, one reproduces the conditions of this definition; cf. Theorem 49. Moreover, at least in the $n = 3$ dimensional vacuum setting, the conditions on initial data on the singularity are such that one locally reproduces the conditions appearing in [29]. If one could localise the existence and uniqueness results of [29], one would also obtain existence and uniqueness of solutions corresponding to the data on the singularity. Moreover, when condition 4 in Definition 10 is void and the data are real analytic, the results of [6] and [22] apply; cf. Subsection 1.5.

Clearly, there are several open problems associated with the topics addressed in this article. First, existence and uniqueness of solutions corresponding to initial data on the singularity should be demonstrated. This should be done using various gauges; cf. the discussion in Subsection 2.2 below. It would also be of interest to generalise the formulations of this article to the case of a non-zero shift vector field. However, it might then be useful to specialise to gauges for which there is a well-posed initial value problem. Moreover, in the context of such gauges, it would be of interest to deduce estimates for higher order derivatives given information concerning a finite number of derivatives; in the context of the asymptotics of solutions to systems of non-linear wave equations, control of a fixed finite number of derivatives usually yields control over all derivatives. The goal would of course be to minimize the number of derivatives involved in Definitions 16 and 17, Theorems 27 and 35 etc. An additional issue is associated with more general matter models. The matter models considered here are isotropic. However, it would be of interest to analyse what happens in the case of anisotropic matter models such as Maxwell’s equations and the Vlasov equation; cf., e.g., [77, 49, 16, 17]. It is also of interest to analyse how the subset of ordinary initial data corresponding to data on the singularity (i.e., conjecturally, the stable
manifold corresponding to the Kasner circle or the Kasner disc) sits in the set of all regular initial data: is this subset a submanifold?

Finally, once a clearer picture of the above issues have been obtained, it is of interest to see if the dynamical systems arguments of, e.g., [18, 59, 23] can be upgraded to the absence of symmetries; cf. Subsection 2.2 below for a further discussion.

1.9. Acknowledgements. This research was funded by the Swedish Research Council, dnr. 2017-03863.

2. Applications

In this section we provide two applications of the notion of initial data on a singularity introduced here.

2.1. A unified perspective on results in the quiescent setting. One advantage of the ideas developed in this article is that they yield a unified perspective on many of the results in the literature. The justification of this statement requires a fairly detailed discussion of a large number of different results, and is, for this reason, quite lengthy. We therefore provide the relevant analysis elsewhere. Nevertheless, we here give one example: we calculate the initial data on the singularity for Bianchi class A solutions to Einstein’s vacuum equations (i.e., the maximal globally hyperbolic development corresponding to left invariant vacuum initial data on 3-dimensional unimodular Lie groups). As discussed in greater detail in [67], the relevant formulation of initial data on the singularity is then the following.

Definition 54. Let \( G \) be a 3-dimensional unimodular Lie group, \( \tilde{h} \) be a left invariant Riemannian metric on \( G \) and \( \mathcal{H} \) be a left invariant (1, 1)-tensor field on \( G \). Then \( (G, \tilde{h}, \mathcal{H}) \) are non-degenerate quiescent Bianchi class A vacuum initial data on the singularity if

1. \( \text{tr} \mathcal{H} = 1 \) and \( \mathcal{H} \) is symmetric with respect to \( \tilde{h} \),
2. \( \text{tr} \mathcal{H}^2 = 1 \) and \( \text{div}_\tilde{h} \mathcal{H} = 0 \),
3. the eigenvalues of \( \mathcal{H} \) are distinct,
4. \( \mathcal{H}_{23} = 0 \) on \( M \), where (if \( p_1 < p_2 < p_3 \) are the eigenvalues of \( \mathcal{H} \), \( \mathcal{H}_A \) is an eigenvector field corresponding to \( p_A \), normalised so that \( |\mathcal{H}_A| = 1 \), and \( \gamma_{BC}^A \) is defined by \( [\mathcal{H}_B, \mathcal{H}_C] = \gamma_{BC}^A \mathcal{H}_A \)), the quantity \( \mathcal{H}_{BC} \) is defined by \( \mathcal{H}_{BC}^A := (\gamma_{BC}^A)^2 \).

As argued in [67], the conditions of Definition 54 yield the existence of an orthonormal (with respect to \( \tilde{h} \)) basis \( \{e_i\} \) of the Lie algebra \( g \) of \( G \), satisfying \( \mathcal{H} e_i = p_i e_i \) (no summation), with \( p_1 < p_2 < p_3 \), as well as \( [e_i, e_j] = \epsilon_{ijk} n_k e_k \) (no summation), where \( \epsilon_{123} = 0 \) and \( \epsilon_{ijk} \) is antisymmetric under permutations of the indices. Here \( p_i \) and \( n_k \) are constants. Moreover, the vanishing/non-vanishing and signs of the \( n_k \) classify the Lie group under consideration: Bianchi type I corresponds to \( n_k = 0 \) for all \( k \); Bianchi type II corresponds to all but one of the \( n_k \) vanishing; Bianchi type VI\( \rho \) corresponds to one of the \( n_k \) vanishing and the remaining ones having opposite sign; Bianchi type VII\( \rho \) corresponds to one of the \( n_k \) vanishing and the remaining ones having the same sign; Bianchi type VIII corresponds to all the \( n_k \) being different from zero and not having the same sign; and Bianchi type IX corresponds to all the \( n_k \) being different from zero and having the same sign. Since \( \mathcal{H}_{23} = 0 \) is, in this setting, equivalent to \( n_1 = 0 \), it is clear that there are no Bianchi type VIII and no Bianchi type IX initial data in the sense of Definition 54. This corresponds to the fact that, disregarding the locally rotationally symmetric solutions (which lead to Cauchy horizons and extendibility of the maximal globally hyperbolic development, cf., e.g., [63] Theorem 24.12, p. 258), all Bianchi type VIII and IX vacuum solutions have oscillatory singularities [59]; i.e., they are not quiescent. However, for the remaining Bianchi class A types, if the solution is not a quotient of a part of Minkowski space, it leads to data on the singularity in the sense of Definition 54 (with the caveat that in some exceptional cases, two \( p_i \)’s could coincide). Moreover, one obtains asymptotics as in Definition 7. We refer the interested reader to [67] for a justification.
of these statements. Finally, in [67], we demonstrate that given data as in Definition 54 there are unique corresponding developments.

To summarise, Definition 54 can be used to parametrise quiescent Bianchi class A vacuum solutions. In the case of Bianchi types VIII and IX, the absence of initial data in the sense of Definition 54 leads to the expectation that solutions should be oscillatory, an expectation which is borne out by [59]. The above is only one very simple example of how previous results can be given a unified understanding by the notion of initial data on the singularity introduced in this article.

2.2. Oscillatory big bang singularities. A second, potential, application is to the understanding of oscillatory big bang singularities. The word “oscillations” here refers to the behaviour of the eigenvalues, say $\ell_A$, of the expansion normalised Weingarten map, say $\mathcal{K}$, along a foliation of a crushing singularity. A singularity is crushing if there is a foliation of a neighbourhood of it such that the mean curvature $\theta$ of the leaves of the foliation diverges to infinity uniformly in the direction of the singularity. In this setting, the expansion normalised Weingarten map is defined by $\mathcal{K} := \bar{K}/\theta$, where $\bar{K}$ denotes the Weingarten map of the leaves of the foliation. In particular, $\text{tr}\mathcal{K} = 1$. Note that $\mathcal{K}$ is symmetric with respect to the induced metric $\bar{g}$, so that the $\ell_A$ are real. Since the sum of the $\ell_A$ equals 1, it is, in 3 + 1-dimensions, convenient to introduce the notation

\begin{equation}
\ell_+ := \frac{3}{2} \left( \ell_2 + \ell_3 - \frac{2}{3} \right) = \frac{3}{2} \left( \frac{1}{3} - \ell_1 \right),
\end{equation}

\begin{equation}
\ell_- := \sqrt{\frac{3}{2}} \left( \ell_2 - \ell_3 \right).
\end{equation}

(57a)

(57b)

Then the $\ell_\pm$ contain all the information concerning the eigenvalues of $\mathcal{K}$, and, in 3 + 1-dimensions the word “oscillations” refers to oscillations in $\ell_\pm$. Moreover, “oscillations” should not here be understood as merely the absence of convergence, but rather a specific type of oscillations illustrated in Figure 1 below.

The origin of the idea that big bang singularities should be oscillatory can be traced back to the physics literature, more particularly the work of Belinski, Khalatnikov and Lifschitz (BKL); cf., e.g., [50, 9, 11]. However, the related ideas of Misner, published around the same time, should also be mentioned, see [55]. The proposal by BKL, often referred to as the BKL conjecture, suggests that the generic dynamics in the direction of a big bang singularity should localise in space, and, along a causal curve, be similar to that of an oscillatory spatially homogeneous solution. In fact, Figure 1 should give a good description of the relevant dynamics. The interested reader is referred to [50, 9, 11] for the original version of these ideas, and to [23, 24, 21, 37, 36] for more recent refinements. Even though the BKL proposal dates back more than 50 years, there are only results concerning oscillatory behaviour in the spatially homogeneous setting; cf., e.g., [59, 71, 60, 35, 18, 8, 49, 15, 25]. This is unsatisfactory, and can be contrasted with a growing number of results concerning the future/past (or both) global non-linear stability of solutions to Einstein’s equations; cf., e.g., [31, 32, 33, 18, 78, 41, 52, 61, 62, 12, 13, 51, 5, 76, 72, 64, 52, 54, 68, 73, 34, 7, 50, 39, 53, 40, 8, 55, 14, 26, 27, 28, 19]. In fact, there are even stability results in the direction of the big bang singularity; cf. [69, 70, 71, 74, 30]. The latter results demonstrate, roughly speaking, that under the circumstances discussed in Subsection 1.5 big bang formation is stable. The settings considered by the authors are the Einstein-scalar field equations, the Einstein-stiff fluid equations and the Einstein vacuum equations in $n$ + 1-dimensions for $n \geq 10$. Moreover, the results concerning the asymptotics include the conclusion that the singularities are quiescent; i.e., that the eigenvalues of $\mathcal{K}$ converge. This may seem to contradict the BKL proposal. However, special matter models and higher dimensions have long been expected to lead to quiescent behaviour. In the physics literature, this expectation appeared in [10, 38]. Two corresponding mathematical results are contained in [4, 22]; cf. Subsection 1.5. Nevertheless, the stability results do not apply in 3 + 1-dimensions unless particular types of matter are present. For this reason, it is of interest to understand oscillatory singularities.
Figure 1. The Kasner circle and the BKL map, using the notation introduced in (57). The circle represents the Kasner solutions; i.e., the geodesically incomplete maximal globally hyperbolic developments corresponding to left invariant vacuum initial data on $\mathbb{R}^3$. Note, in particular, that the solutions are parametrised by the eigenvalues of the expansion normalised Weingarten map. Since these solutions are invariant under permutations of the eigenvalues, there is a sixfold symmetry in this picture, meaning that only one sixth of the circle represents geometrically distinct solutions. In particular, it is, e.g., sufficient to focus on the segment between $Q_1$ (the point antipodal to $T_1$) and $T_3$. The points $T_i$ represent the flat Kasner solutions (the corresponding Riemann curvature tensor is identically zero) and the $T_i$ together with their antipodal points, denoted $Q_i$, represent the locally rotationally symmetric Kasner solutions. Given a point $p$ on the circle, the BKL map applied to $p$, denoted $\kappa(p)$, is obtained as follows: take the corner of the triangle closest to $p$; draw a straight line from this corner to $p$; continue the straight line to the next intersection with the circle; and define $\kappa(p)$ to be this second intersection. In the figure, we illustrate four iterations of $\kappa$ on a specific point. The chaoticity of the BKL map follows from, e.g., [8, Section 8, p. 22].

Due to the above observations, it is clear that it is important to prove results in the spatially inhomogeneous and oscillatory setting. Given the existing results, it is therefore natural to ask if ideas similar to those used to prove stability, in particular the arguments of [30], can be used. The answer to this question is most probably no, and the main reason is that the background solutions considered in [30] are orbitally stable, whereas the model dynamics in the oscillatory setting are chaotic; cf. Figure 1. In particular, the model behaviour in the oscillatory setting is by its very nature unstable. The arguments required to make progress on analysing this setting can thus be expected to be quite different. Next, since there are, so far, only results in the spatially homogeneous setting, it is natural to ask if the corresponding arguments can be used. Unfortunately, the ideas used in the first results in this setting are based on monotonicity principles and contradiction arguments; cf., e.g., [59, 77, 60, 35, 15]. In particular, global existence is guaranteed by soft methods. Moreover, contradiction arguments (combined with, e.g., the monotonicity principle) are used to derive information concerning the asymptotics. In the spatially inhomogeneous setting, it is unlikely to be possible to derive global existence without at the same time quantitatively controlling the asymptotic behaviour of solutions. Moreover, to the extent that there are monotonic quantities, these quantities can not be expected to yield much information
concerning the asymptotics. To conclude, the ideas of \cite{59,74,60,35,15} are unlikely to be useful when trying to obtain results in the spatially inhomogeneous setting. On the other hand, more recently, results based on estimates have appeared; cf., e.g., \cite{48,8,49,25}. The idea of these articles is to take an orbit of the model system, cf. Figure 1 and then to prove that there is a stable manifold of solutions to the actual equations converging to this orbit. In \cite{48,8,49}, the union of the corresponding stable manifolds is conjectured to be non-generic. However, in \cite{25}, the author treats general enough orbits of the model system that the union of the corresponding stable manifolds has positive Lebesgue measure in the relevant set of initial data.

In order to understand how the ideas of \cite{48,8,49,25} could, potentially, be used in the spatially inhomogeneous setting, it is convenient to focus on a specific situation. In what follows, we therefore discuss vacuum Bianchi type IX solutions, the goal being to give a rough idea of the arguments in \cite{48,25}. Note, to begin with, that generic vacuum Bianchi type IX solutions converge therefore discuss vacuum Bianchi type IX solutions, the goal being to give a rough idea of the inhomogeneous setting, it is convenient to focus on a specific situation. In what follows, we

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corresponding to the eigenvalue \( p \) coincide). The reason for this is that the only degenerate points on the Kasner circle are the \( T_i \)’s and their antipodal points; cf. Figure 1. On the other hand, orbits of the model system which include \( T_i \) (or its antipodal point) terminate at \( T_i \). In the case of Bianchi type VIII and IX, it can be demonstrated that there are no solutions to the actual equations that terminate at the point \( T_i \). However, there are solutions whose \((\ell_+ \ell_-)\)-coordinates terminate on a \( T_i \). On the other hand, those solutions are locally rotationally symmetric (cf. [63 Proposition 22.10, p. 237]) and have a Cauchy horizon through which the solution can be extended (cf. [63 Theorem 24.12, p. 258]). In particular, these solutions are non-generic and do not reflect the behaviour of generic solutions. For these reasons, the stable manifold associated with degenerate points on the Kasner circle is not expected to be of any greater interest.

The proofs in [48, 25] are highly non-trivial. Nevertheless, the arguments are substantially simplified not only by the fact that the authors consider spatially homogeneous solutions, but by the fact that they consider Bianchi class A solutions (maximal globally hyperbolic developments corresponding to left invariant initial data on 3-dimensional unimodular Lie groups). The reason for the simplification is that in the unimodular setting, the stable and unstable manifolds associated with the Kasner circle correspond to the vanishing/sign of basic variables that, at the same time, classify the Lie group under consideration; cf. [63 Section 3.4] for a more detailed discussion. Already the Bianchi class B (i.e., the non-unimodular) setting is more complicated, but it is reasonable to expect the spatially inhomogeneous setting to represent a different level of difficulty. This is partly due to the problem of gauge invariance; there is a large number of choices of gauge (including choices of foliation). Moreover, instead of the simple dichotomy of being on the stable manifold or off (as in the Bianchi class A setting), spatial variations can be expected to give rise to complications such as spikes; cf. [58] for a discussion of this topic in the context of \( \mathbb{T}^3 \)-Gowdy vacuum spacetimes. Causal localisation might also be necessary in order to deal with the substantial spatial variations that are to be expected. However, such a localisation would limit the allowed gauge choices. Since the causal structure is not known a priori, causal localisation also makes it necessary to keep track of the causal structure “along an orbit”. And so on. Finally, it is not clear that it would be possible to carry out ideas analogous to [48, 25] in the general spatially inhomogeneous setting. Nevertheless, at present, this approach seems to be the most promising.

If one wants to generalise the ideas of [48, 25] to the spatially inhomogeneous setting, the natural first step is to identify the stable manifold corresponding to the Kasner circle in the absence of symmetries. However, as pointed out in Subsection 1.8, it is reasonable to conjecture that Definition 2 parametrises this stable manifold.

3. RELATING DIFFERENT NOTIONS OF INITIAL DATA, THE 3-DIMENSIONAL VACUUM SETTING

The purpose of the present section is to relate the two notions of initial data on the singularity given by Definitions 1 and 2. We start by, given initial data in the sense of Definition 2, constructing appropriate local coordinates.

**Lemma 55.** Let \( \tilde{M} \) be a 3-dimensional manifold and \( \mathcal{X} \) be a smooth \((1,1)\)-tensor field on \( \tilde{M} \). Assume \( \mathcal{X} \) to have distinct real eigenvalues \( p_1 < p_2 < p_3 \). Let \( \tilde{\mathcal{F}}_A \) be an eigenvector field of \( \mathcal{X} \) corresponding to the eigenvalue \( p_A \) and let \( \gamma^A_{BC} \) be defined by \([\tilde{\mathcal{F}}_B, \tilde{\mathcal{F}}_C] = \gamma^A_{BC} \tilde{\mathcal{F}}_A \). Assume that \( \gamma^1_{23} = 0 \). Then, if \( p \in \tilde{M} \), there are local coordinates \((V, \tilde{x})\) with \( p \in V \), such that, after a renormalisation of the \( \tilde{\mathcal{F}}_A \), if necessary,

\[
\begin{align*}
\mathcal{F}_3 &= \partial_3, \\
\mathcal{F}_2 &= \partial_2 + \mathcal{F}_2^3 \partial_3, \\
\mathcal{F}_1 &= \partial_1 + \mathcal{F}_1^2 \partial_2 + \mathcal{F}_1^3 \partial_3
\end{align*}
\]

for a suitable choice of smooth functions \( \mathcal{F}_1^2, \mathcal{F}_2^3, \mathcal{F}_1^3 \). Moreover, if \( \{\tilde{\mathcal{F}}_A\} \) is the frame dual to \( \{\mathcal{F}_A\} \), then

\[
\begin{align*}
\mathcal{F}^1 &= d\tilde{x}^1, \\
\mathcal{F}^2 &= d\tilde{x}^2 - \mathcal{F}_1^2 d\tilde{x}^3, \\
\mathcal{F}^3 &= d\tilde{x}^3 - \mathcal{F}_2^3 d\tilde{x}^2 + (\mathcal{F}_1^2 \mathcal{F}_2 - \mathcal{F}_1^3) d\tilde{x}^1.
\end{align*}
\]
Define $\mathcal{K}_i^j$ by the condition that $\mathcal{K} \partial_i = \mathcal{K}_i^j \partial_j$. Then $\mathcal{K}_i^j = 0$ if $j < i$; $\mathcal{K}_i^i = p_i$ (no summation); and

\begin{align}
(60a) & \quad \mathcal{F}_1^2 = \frac{1}{p_1 - p_2} \mathcal{K}_1^2, \\
(60b) & \quad \mathcal{F}_2^3 = \frac{1}{p_2 - p_3} \mathcal{K}_2^3, \\
(60c) & \quad \mathcal{F}_1^3 = \frac{1}{p_1 - p_3} \mathcal{K}_1^3 + \frac{1}{(p_1 - p_3)(p_1 - p_2)} \mathcal{K}_1^2 \mathcal{K}_2^3.
\end{align}

**Remark 56.** In [29, Theorem 1.1, pp. 2-3], a quantity $\kappa_i^j$ appears. Due to the differences in convention, $\mathcal{K}_i^j = -\kappa_i^j$.

**Proof.** To begin with, there are local coordinates around $p$ such that $\mathcal{F}_3 = \partial_3$; cf., e.g., [37, Theorem 9.22, p. 220]. By assumption

$$[\partial_3, \mathcal{F}_2] = A \partial_3 + B \mathcal{F}_2.$$ 

In order to find a vector field which commutes with $\partial_3$, and which is, at the same time, in the span of $\partial_3$ and $\mathcal{F}_2$, let $X_2 = \alpha \mathcal{F}_2 + \beta \partial_3$. Then

$$[\partial_3, X_2] = (\partial_3 \alpha) \mathcal{F}_2 + \alpha [\partial_3, \mathcal{F}_2] + (\partial_3 \beta) \partial_3 = (\partial_3 \alpha + B \alpha) \mathcal{F}_2 + (\partial_3 \beta + A \alpha) \partial_3.$$ 

Assume that $p$ corresponds to the origin in the local coordinates, and fix $\alpha = 1$ on the plane defined by $x^3 = 0$. Then we can calculate $\alpha$ by solving the equation $\partial_3 \alpha + B \alpha = 0$. Moreover, $\alpha$ will be smooth and strictly positive in a neighbourhood of the origin. Next, we solve for $\beta$ by choosing $\beta = 0$ on the plane $x^3 = 0$ and integrating $\partial_3 \beta + A \alpha = 0$. Doing so yields a vector field $X_2$ with the property that the span of $X_2$ and $\partial_3$ equals the span of $\mathcal{F}_2$ and $\mathcal{F}_3$ in a neighbourhood of $p$. Moreover, $\partial_3$ and $X_2$ commute. Due to [47, Theorem 9.46, p. 234], we conclude that there are local coordinates in a neighbourhood of $p$ such that $\mathcal{F}_3 = \partial_3$ and $X_2 = \partial_2$. Since $\mathcal{F}_2$ is in the span of $\mathcal{F}_3$ and $X_2$, it is clear that $\mathcal{F}_2 = a \partial_2 + b \partial_3$. Since $\mathcal{F}_2$ and $\mathcal{F}_3$ are linearly independent, it is clear that $a$ is never allowed to vanish. In particular, after renormalising the eigenvector field $\mathcal{F}_2$, if necessary, we can assume that $\mathcal{F}_2 = \partial_2 + \mathcal{F}_2^3 \partial_3$. That $\mathcal{F}_1$ can be assumed to take the form $\mathcal{F}_1 = \partial_1 + \mathcal{F}_1^2 \partial_2 + \mathcal{F}_1^3 \partial_3$ follows by a similar argument.

Finally, (59) follows from straightforward calculations and (60) and the statements concerning $\mathcal{K}_i^j$ follow from the fact that $\mathcal{F}_A$ is an eigenvector field of $\mathcal{K}$ with eigenvalue $p_A$. \qed

### 3.1. Relating the limits.

In Proposition 3 the goal is to, locally, associate data in the sense of Definition 1 with initial data in the sense of Definition 2. In the present section we provide the intuition behind the definition (12). We do so by assuming that we have a development in the sense of Definition 7. We then prove that this development can locally be written in the form (14). Finally, we calculate the limit (5).

**Lemma 57.** Let $(\tilde{M}, \tilde{h}, \tilde{\mathcal{K}})$ be non-degenerate quiescent vacuum initial data on the singularity in the sense of Definition 2. Assume that there is a locally Gaussian development $(M, g)$ of the initial data in the sense of Definition 7. Fix a $p \in \tilde{M}$ and let $(V, \tilde{x})$ be local coordinates with the properties stated in Lemma 55. Then, on $V \times (0, t_+)$, the development (13) can be written

$$\Psi' = -dt \otimes dt + \sum_{i,j=1}^{3} a_{ij} t^{2p_{\max(i,j)}} d\tilde{x}^i \otimes d\tilde{x}^j,$$

using the notation of Definition 7. Moreover, (5) is satisfied, where $c_{ii} = \tilde{h}(\tilde{x}_i, \tilde{x}_i)$ (no summation) and the $\tilde{x}_i$ are given in the statement of Lemma 55: the $c_{i,j}$, $i < j$, are given by (12), where the sub- and superscripts refer to the coordinate frame associated with $(V, \tilde{x})$; and $c_{ij} = c_{ji}$. Moreover, if $W$ has compact closure contained in $V$, there is an $\eta > 0$, a $t_0 \in (0, t_+)$ and, for every $l \in \mathbb{N}$, a $C_l$ such that

$$\|a_{ij}(\cdot, t) - c_{ij}|_{C^l(W)} \leq C_l t^{2\eta}$$

on $(0, t_0)$. 

Define the $p_i$ and $c_{ij}$ as above and $\kappa_i^j = -\mathcal{H}_i^j$. Then the $p_i$ and $c_{ij}$ are smooth real valued functions on $V$ and the $p_i$ satisfy (4). Moreover, $\kappa_i^j = -p_i$ (no summation); $\kappa_i^j = 0$ if $l < i$; and (3) is satisfied.

**Remark 58.** Arguments similar to the proof of the lemma provide a proof of all the statements of Proposition 6 except for the statement that (2) holds. They also demonstrate that $\kappa_i^j = -\mathcal{H}_i^j$.

**Proof.** Let $\mathcal{X}_A$ be the vector fields defined by (58). Note that $\mathcal{X}_A \mathcal{X}_B = p_A \mathcal{X}_A$ (no summation). Since $\mathcal{X}$ is symmetric with respect to $\mathcal{h}$, the $\mathcal{X}_A$ are orthogonal with respect to $\mathcal{h}$. In particular, there are thus smooth functions $\beta_A > 0$ such that

\[
\mathcal{h} = \sum_A \beta_A^2 \mathcal{X}_A \otimes \mathcal{X}_A.
\]

Note that $\beta_A^2 = \mathcal{h}(\mathcal{X}_A, \mathcal{X}_A)$ (no summation). Combining this observation with (15b) yields the conclusion that, for $A \neq B$, $b_{AB} \to 0$ as $t \to 0^+$, where the convergence is uniform on compact subsets of $V$. Moreover, $b_{AA} \to \beta_A^2$ as $t \to 0^+$ (no summation), where the convergence is uniform on compact subsets of $V$.

Next, we wish to rewrite the metric in terms of the coordinate basis $dx^i \otimes dx^j$, where $(V, \tilde{x})$ are the coordinates constructed in Lemma 55. Note, to this end, that

\[
\sum_{A,B} b_{AB} t^{2\max(A,B)} \mathcal{X}_A \otimes \mathcal{X}_B = \sum_{i,j} a_{ij} t^{2\max(i,j)} dx^i \otimes dx^j,
\]

where

\[
a_{ij} = \sum_{A,B} b_{AB} t^{2\max(A,B)} \mathcal{X}_A \otimes \mathcal{X}_B = \sum_{i,j} a_{ij} t^{2\max(i,j)} dx^i \otimes dx^j.
\]

Here, due to (59), $\mathcal{X}_i^A = 0$ if $A < i$, $\mathcal{X}_i^i = 1$ if $A = i$ and

\[
\mathcal{X}_1^2 = -\mathcal{X}_2^2, \quad \mathcal{X}_1^3 = -\mathcal{X}_2^3, \quad \mathcal{X}_1^1 = \mathcal{X}_2^1 - \mathcal{X}_1^3.
\]

Consider $a_{ii}$ (no summation). If $A > i$ or $B > i$, then the corresponding term in (63) converges to zero, since the $b_{AB}$ converge and $2\max(A,B) - 2\max(i,j) > 0$ in that case. The only term we need to take into consideration is thus the one with $i = j$ and $A = B = i$. This term reads $b_{ii}$ (no summation) and, by the above, it converges to $\beta_i^2$. To conclude, $a_{ii}$ converges to $\beta_i^2$, and in order for the notation to be consistent with that of [29], we introduce $c_{ii} = \beta_i^2$ (no summation). In particular $c_{ii}$ is thus determined by $\mathcal{h}$ and our choice of frame; cf. (62).

Next, consider $a_{12}$. If $A$ or $B$ equals 3, then the corresponding term in (63) converges to zero. If $A = B = 1$, the corresponding term equals zero, since $\mathcal{X}_1^1 = 0$. The only possibilities for $(A, B)$ that remain to be considered are $(1, 2)$, $(2, 1)$ and $(2, 2)$. However, since $b_{12}$ converges to zero, it is sufficient to consider

\[
b_{22} \mathcal{X}_2^2 \mathcal{X}_2^2 \to -\mathcal{X}_2^2 c_{22},
\]

where we used the fact that $b_{22}$ converges to $c_{22}$, the fact that $\mathcal{X}_2^2 = 1$ and the fact that $\mathcal{X}_1^1 = -\mathcal{X}_2^2$. Again, to be consistent with [29], we use the notation $c_{12} := -\mathcal{X}_1^2 c_{22}$ (and note that $a_{12}$ converges to $c_{12}$). Combining this relation with (60b) yields

\[
\mathcal{X}_2^2 = (p_2 - p_1) c_{12} c_{22}.
\]

Letting $\kappa_i^j = -\mathcal{H}_i^j$, (64) contains the same information as the first equality in (3). Next, consider $a_{23}$. Then $B$ has to equal 3. Moreover, if $A < 3$, the corresponding term converges to zero since $b_{AB} \to 0$ for $A \neq B$. The only term we need to consider is thus

\[
b_{33} \mathcal{X}_3^3 \mathcal{X}_3^3 = -\mathcal{X}_3^3 b_{33} \to -\mathcal{X}_3^3 c_{33}.
\]

In particular, $a_{23}$ converges to $c_{23} := -\mathcal{X}_2^3 c_{33}$. Keeping (60b) in mind, we obtain

\[
\mathcal{X}_3^3 = (p_3 - p_2) c_{23} c_{33},
\]

cf. the second equality in (3). Finally, consider $a_{13}$. For the same reasons as above, the only term we have to take into consideration is

\[
b_{33} \mathcal{X}_1^3 \mathcal{X}_3^3 = b_{33} (\mathcal{X}_1^2 \mathcal{X}_3^3 - \mathcal{X}_1^3) \to c_{33} (\mathcal{X}_1^2 \mathcal{X}_3^3 - \mathcal{X}_1^3) =: c_{13}.
\]
Appealing to (60) as well as (64) and (65) yields the conclusion that
\[ \mathcal{K}_i^3 = (p_3 - p_1) \frac{c_{13}}{c_{33}} - (p_2 - p_1) \frac{c_{12} c_{23}}{c_{22} c_{33}}; \]
cf. the third equality in (3). Due to \( c_{ij} = \beta^2 \), (64), (65), and (66), it is clear that if the \( c_{ij} \) are defined by (12) for \( i < j \), then (5) holds. In fact, if \( W \) has compact closure contained in \( V \), (61) holds.

The final statements of the lemma follow from the fact that \( \text{tr} \mathcal{K} = \text{tr} \mathcal{K}^2 = 1 \), the conclusions of Lemma 55 and the above observations.

3.2. The asymptotic momentum constraint. Due to Remark 58 what remains (in order to prove Proposition 6) is to prove that (2) holds.

Lemma 59. Let \((\tilde{M}, \tilde{h})\) be a smooth 3-dimensional Riemannian manifold and \( \mathcal{K} \) be a smooth \((1, 1)\)-tensor field on \( \tilde{M} \) which is symmetric with respect to \( \tilde{h} \). Assume \( \mathcal{K} \) to have distinct eigenvalues \( p_1 < p_2 < p_3 \). Assume, moreover, that there are local coordinates \((V, \tilde{x})\) such that for \( A = 1, 2, 3 \), there are eigenvector fields \( \tilde{\mathcal{F}}_A \), \( A = 1, 2, 3 \), of \( \mathcal{K} \) of the form (58) for a suitable choice of smooth functions \( \tilde{\mathcal{F}}^1, \tilde{\mathcal{F}}^2, \tilde{\mathcal{F}}^3 \). Assume also that \( \mathcal{K} \tilde{\mathcal{F}}_A = p_A \tilde{\mathcal{F}}_A \) (no summation) for \( A = 1, 2, 3 \). Define \( c_{ij} \) as in the statement of Proposition 6. Then (2) (where \( \kappa_i^j \) is defined in terms of \( c_{ij} \) and \( p_i \)) as described in Definition 7 is equivalent to \( \text{div}_h \mathcal{K} = 0 \).

Proof. Define \( \{\tilde{\mathcal{F}}^A\} \) via (59). Since \( \mathcal{K} \) is symmetric with respect to \( \tilde{h} \), it follows that \( \{\tilde{\mathcal{F}}_A\} \) is an orthogonal frame with respect to \( \tilde{h} \). Thus
\[ \tilde{h} = \sum_A c_{AA} \tilde{\mathcal{F}}^A \otimes \tilde{\mathcal{F}}^A, \]
where \( c_{AA} := \tilde{h}(\tilde{\mathcal{F}}_A, \tilde{\mathcal{F}}_A) \) (no summation); cf. Proposition 6. Next, note that
\[
\tilde{D}_A \mathcal{K}^A_B = \tilde{\mathcal{F}}_A(\mathcal{K}^A_B) - (\tilde{D}_A \tilde{\mathcal{F}}_A)(\mathcal{K} \tilde{\mathcal{F}}_B) - \tilde{\mathcal{F}}^A(\mathcal{K} \tilde{D}_A \tilde{\mathcal{F}}_B)
= \tilde{\mathcal{F}}_A(p_B \beta^A) + \tilde{\mathcal{F}}^A(\tilde{D}_A \mathcal{F}_C) \mathcal{F}^C - \tilde{\mathcal{F}}^A(\mathcal{K} \tilde{D}_A \mathcal{F}_B)
= \tilde{D}_B(p_B) + p_B \mathcal{F}^A(\tilde{D}_A \mathcal{F}_B) - \tilde{\mathcal{F}}^A(\mathcal{K} \tilde{D}_A \tilde{\mathcal{F}}_B)
\]
(no summation on \( B \)), where \( \tilde{D} \) denotes the Levi-Civita connection induced by \( \tilde{h} \). Introducing the notation \( \Gamma^C_{AB} \) by the requirement that
\[ \tilde{D}_A \mathcal{F}_B = \Gamma^C_{AB} \mathcal{F}_C; \]
this equality can be written (again, no summation on \( B \))
\[ \tilde{D}_A \mathcal{K}^A_B = \mathcal{F}_B(p_B) + \sum_A(p_B - p_A) \Gamma^A_{AB}. \]
What remains is thus to compute (no summation on \( A \))
\[ \Gamma^A_{AB} = \frac{1}{c_{AA}} \tilde{h}(\tilde{D}_A \tilde{\mathcal{F}}_B, \tilde{\mathcal{F}}_A) = \frac{1}{c_{AA}} \tilde{h}([\tilde{\mathcal{F}}_A, \tilde{\mathcal{F}}_B] + \tilde{D}_B \tilde{\mathcal{F}}_A, \tilde{\mathcal{F}}_A) = \frac{1}{2} \mathcal{F}_B(\ln c_{AA}) + \gamma^A_{AB}. \]
where \( \gamma^A_{BC} \) is defined by (11). Next, we need to calculate \( \gamma^A_{AB} \) (no summation). Keeping (58) and (59) in mind, it can be computed that the only non-zero \( \gamma^A_{AB} \) (no summation) are
\[ \gamma^2_{21} = \partial_2 \mathcal{F}_1^2 + \mathcal{F}_3 \partial_3 \mathcal{F}_1^2, \]
\[ \gamma^3_{31} = \partial_3 \mathcal{F}_1^3 - \mathcal{F}_2 \partial_2 \mathcal{F}_1^3, \]
\[ \gamma^3_{32} = \partial_3 \mathcal{F}_2^3. \]
Combining the above observations, we can compute that
\[ 2 \mathcal{D}_A \mathcal{K}^A_B = 2 \partial_B p_3 + \sum_A(p_3 - p_A) \partial_3 \ln c_{AA}. \]
In particular, \(-2\tilde{D}_A \mathcal{X}_3^A\) equals the left hand side of (2) in case \(i = 3\). In particular, \(\tilde{D}_A \mathcal{X}_3^A = 0\) is equivalent to (2) in case \(i = 3\). Next, compute
\[
2\tilde{D}_A \mathcal{X}_2^A = 2\tilde{X}_2(p_2) + \sum_A (p_2 - p_A) \tilde{X}_2^2 \left(\ln c_{AA}\right) + 2(p_2 - p_3) \partial_3 \tilde{X}_2^3
\]
\[
= 2\partial_1 p_2 + 2\tilde{X}_2^2 \partial_3 p_2 + \sum_A (p_2 - p_A) \partial_3 \left(\ln c_{AA}\right)
+ \tilde{X}_2^3 \sum_A (p_2 - p_A) \partial_3 \left(\ln c_{AA}\right) + 2(p_2 - p_3) \partial_3 \tilde{X}_2^3.
\]

At this stage, it is of interest to note that
\[
2\tilde{X}_2^3 \partial_3 p_2 + \tilde{X}_2^3 \sum_A (p_2 - p_A) \partial_3 \left(\ln c_{AA}\right) + 2(p_2 - p_3) \partial_3 \tilde{X}_2^3
\]
\[
= \tilde{X}_2^3 \left[2\partial_3 p_3 + \sum_A (p_3 - p_A) \partial_3 \left(\ln c_{AA}\right)\right] + 2\tilde{X}_2^3 \partial_3 (p_2 - p_3)
+ \tilde{X}_2^3 \sum_A (p_2 - p_A) \partial_3 \left(\ln c_{AA}\right) + 2(p_2 - p_3) \partial_3 \tilde{X}_2^3
\]
\[
= 2\tilde{X}_2^3 \tilde{D}_A \mathcal{X}_3^A + \mathcal{X}_2^3 \sum_A \partial_3 \left(\ln c_{AA}\right) + 2\partial_3 \mathcal{X}_2^3,
\]
where we appealed to (60) in the last step. Summing up,
\[
2\tilde{D}_A \mathcal{X}_2^A - 2\tilde{X}_2^3 \tilde{D}_A \mathcal{X}_3^A
= \sum_A (p_2 - p_A) \partial_3 \left(\ln c_{AA}\right) + 2\partial_3 p_2 + 2\partial_3 \mathcal{X}_2^3 + \mathcal{X}_2^3 \sum_A \partial_3 \left(\ln c_{AA}\right).
\]

In particular, minus the right hand side of this equation equals the left hand side of (2) in case \(i = 2\).

Finally, let us consider
\[
2\tilde{D}_A \mathcal{X}_1^A - 2\tilde{X}_1^2 \tilde{D}_A \mathcal{X}_2^A - 2\tilde{X}_1^2 \tilde{D}_A \mathcal{X}_3^A + 2\tilde{X}_1^2 \tilde{D}_A \mathcal{X}_3^A
\]
\[
= 2\tilde{X}_1(p_1) + \sum_A (p_1 - p_A) \tilde{X}_1^2 \left(\ln c_{AA}\right) + 2(p_1 - p_2) \partial_2 \tilde{X}_1^2
+ 2(p_1 - p_3) \partial_3 \tilde{X}_1^2 - 2\tilde{X}_1^2 \tilde{X}_2(p_2) - \sum_A (p_2 - p_A) \tilde{X}_1^2 \tilde{X}_2^2 \left(\ln c_{AA}\right) - 2(p_2 - p_3) \tilde{X}_1^2 \partial_3 \tilde{X}_2^3
- 2(\tilde{X}_2^3 - \tilde{X}_2^3 \tilde{X}_2^3) \partial_3 p_3 - \sum_A (p_3 - p_A) (\tilde{X}_2^3 - \tilde{X}_2^3 \tilde{X}_2^3) \partial_3 \left(\ln c_{AA}\right).
\]

Collecting the terms on the right hand side involving sums over \(A\) yields
\[
\sum_A (p_1 - p_A) \partial_1 \left(\ln c_{AA}\right) + \mathcal{X}_1^2 \sum_A \partial_2 \left(\ln c_{AA}\right) + \mathcal{X}_1^3 \sum_A \partial_3 \left(\ln c_{AA}\right).
\]
The remaining terms sum up to
\[
2\partial_1 p_1 + 2\partial_2 \mathcal{X}_2^2 + 2\partial_3 \mathcal{X}_2^3.
\]

To conclude
\[
2\tilde{D}_A \mathcal{X}_1^A - 2\tilde{X}_1^2 \tilde{D}_A \mathcal{X}_2^A - 2\tilde{X}_1^2 \tilde{D}_A \mathcal{X}_3^A + 2\tilde{X}_1^2 \tilde{D}_A \mathcal{X}_3^A
\]
\[
= \sum_A (p_1 - p_A) \partial_1 \left(\ln c_{AA}\right) + 2\partial_1 \mathcal{X}_1^A + \mathcal{X}_1^2 \sum_A \partial_2 \left(\ln c_{AA}\right) + \mathcal{X}_1^3 \sum_A \partial_3 \left(\ln c_{AA}\right).
\]

Moreover, the right hand side equals minus the left hand side of (2) in case \(i = 1\). The statement follows.

Finally, we prove Proposition 6.

**Proposition 6.** The conclusion of Proposition 6 is an immediate consequence of Remark 58 and Lemma 59.

3.3. **Deriving geometric asymptotics.** In Subsection 3.2 we verified that, locally, initial data in the sense of Definition 2 are initial data in the sense of Definition 1. Next, we wish to verify that if we have a development in the sense of 29, we obtain a development in the sense of Definition 7.

**Proposition 60.** Let \((\bar{M}, \bar{h})\) be a smooth 3-dimensional Riemannian manifold and \(\mathcal{X}\) be a smooth \((1,1)\)-tensor field on \(M\) which is symmetric with respect to \(\bar{h}\). Assume \(\mathcal{X}\) to have distinct eigenvalues \(p_1 < p_2 < p_3\). Assume that there are local coordinates \((V, \bar{x})\) such that for \(A = 1, 2, 3\), there are eigenvector fields \(\bar{X}_A\), \(A = 1, 2, 3\), of \(\mathcal{X}\) of the form \(\mathcal{X}_A = p_A \bar{X}_A\) for a suitable choice of smooth functions \(\bar{X}_1^A, \bar{X}_2^A, \bar{X}_3^A\). Assume also that \(\mathcal{X}_A \bar{X}_A = p_A \bar{X}_A\) (no summation) for \(A = 1, 2, 3\). Define
c_{ij}$ as in the statement of Proposition 6 and $\kappa_{ij}$ as in Definition 1. Assume, finally, that there is a solution to Einstein’s vacuum equations of the form (4) on $V \times (0, t_0)$ for some $t_0 > 0$, where $dx^i$ is replaced by $dx^i$, such that

\begin{align}
(67) & \quad \|a_{ij}(\cdot, t) - c_{ij}\|_{C^r(W)} \leq C_l t^\varepsilon, \\
(68) & \quad \|t\kappa_{ij}(\cdot, t) + \kappa_{ij}\|_{C^r(W)} \leq C_l t^\varepsilon
\end{align}
onumber

on $(0, t_0)$ for some constant $C_l$, where $l \in \mathbb{N}$, $C_l \in \mathbb{R}$ and $W \subset M$ is such that it has compact closure contained in $V$. Then the development (4) can be represented by the right hand side of (13) on $V \times (0, t_0)$, where $\mathcal{Y}^A$ is replaced by $\mathcal{Y}^A$. Moreover, defining $h$ by (14), where $\mathcal{Y}^A$ is replaced by $\mathcal{Y}^A$, there is an $\eta > 0$ and a constant $C_l$ such that

\begin{align}
(69) & \quad \|K(\cdot, t) - \mathcal{K}\|_{C^r(W)} \leq C_l t^\eta, \\
(70) & \quad \|\bar{h}(\cdot, t) - \bar{h}\|_{C^r(W)} \leq C_l t^\eta,
\end{align}

hold on $(0, t_0)$.

Proof. First note that the definitions ensure that $\kappa_{ij} = -\mathcal{K}_{ij}$. Due to (68), it is clear that $t\theta$ converges to 1 in $C^l$ at the rate $t^\varepsilon$. Combining this information with (68) again yields the conclusion that (69) holds. It remains to relate $h$ to $\bar{h}$. Note, to this end, that

\[ \bar{h} = \sum A h(\mathcal{F}_A, \mathcal{F}_A) \mathcal{Y}^A \otimes \mathcal{Y}^A. \]

Note, moreover, that $dx^i = \mathcal{F}_A \mathcal{Y}^A$, where the $\mathcal{F}_A$ are defined by (58); $\mathcal{F}_A = 0$ if $i < A$; and $\mathcal{F}_A = 1$ if $A = i$. Thus the metric can be represented by the right hand side of (13), with $\mathcal{Y}^A$ replaced by $\mathcal{Y}^A$, where

\begin{align}
(71) & \quad b_{AB} = \sum_{i,j} a_{ij} \varepsilon^{2p_{\text{max}}(i,j) - 2p_{\text{max}}(A,B)} \mathcal{F}_A \mathcal{F}_B.
\end{align}

Consider, to begin with, $b_{AA}$ (no summation). If $i$ or $j$ in (71) is strictly larger than $A$, then the corresponding term in (71) converges to zero. Moreover, we have to have $i \geq A$ and $j \geq A$ since $\mathcal{F}_A = 0$ if $i < A$. To conclude, the only term in (71) which contributes to the limit is the one with $i = j = A$. Since $\mathcal{F}_A = 1$ if $A = i$, this term equals $a_{AA} \rightarrow c_{AA} = \bar{h}(\mathcal{F}_A, \mathcal{F}_A)$ (no summation). In other words, $b_{AA} \rightarrow \bar{h}(\mathcal{F}_A, \mathcal{F}_A)$ (no summation), as desired. Next, note that

\[ b_{23} = a_{33} \mathcal{F}_2 \mathcal{F}_3 + a_{23} \mathcal{F}_2 \mathcal{F}_3 = a_{33} \mathcal{F}_2 \mathcal{F}_3 + a_{23} \rightarrow c_{33} \mathcal{F}_2 \mathcal{F}_3 + c_{23} = 0, \]

by the definition of $c_{23}$. Next consider $b_{12}$. In this case the only terms in (71) that contribute to the limit are the ones with $j = 2$ and $i \in \{1, 2\}$:

\[ a_{12} \mathcal{F}_1 \mathcal{F}_2 + a_{22} \mathcal{F}_1 \mathcal{F}_2 = a_{12} + a_{22} \mathcal{F}_1 \mathcal{F}_2 \rightarrow c_{12} + c_{22} \mathcal{F}_1 \mathcal{F}_2 = 0, \]

by the definition of $c_{12}$. Finally, consider $b_{13}$. In this case, the only terms in (71) that contribute to the limit are the ones with $j = 3$ and $i \in \{1, 2, 3\}$:

\[ a_{13} + a_{23} \mathcal{F}_1 \mathcal{F}_3 + a_{33} \mathcal{F}_1 \mathcal{F}_3 \rightarrow c_{13} + c_{23} \mathcal{F}_1 \mathcal{F}_3 + c_{33} \mathcal{F}_1 \mathcal{F}_3 = 0, \]

by the definition of $c_{13}$ and $c_{23}$. The lemma follows. \hfill \Box

4. Existence of developments

The purpose of the present section is to prove Propositions 18 and 21 as well as Theorem 22.

Proposition 18. In what follows, we refer to (6), which covers the case $n = 3$. However, the equations for the velocity dominated solutions are the same in higher dimensions; cf. (22). Nevertheless, in all the verifications to follow, Remark 19 should be kept in mind. Define $(0g, 0^k, 0^\phi)$ by (19). Then $0^k A = -p_A^{-1} \delta^k_B$ (no summation), where the indices are raised with $0^g$, and the indices refer to the frame $\{\mathcal{F}_A\}$ and coframe $\{\mathcal{Y}^A\}$ appearing in the statement of the proposition. One
particular consequence of this equality is that $\text{tr}^0 k = - t^{-1}$. Note also that $\partial_t^0 g_{AB} = -2^0 k_{AB}$; i.e., [6] (11a), p. 483] is satisfied. Next we define, in analogy with [6] (12), p. 483],
\begin{align}
(72a) \quad 0 \rho & := \frac{1}{2} (\partial_t^0 \phi)^2, \\
(72b) \quad 0 j_A & := - (\partial_t^0 \phi) \cdot \mathcal{X}_A^0 (\phi), \\
(72c) \quad 0 S_{AB} & := \frac{1}{2} (\partial_t^0 \phi)^2 \cdot (0 g_{AB}).
\end{align}

With these definitions, the matter terms appearing on the right hand side of [6] (11b), p. 483] vanish. Due to this observation, it can be verified that [6] (11b), p. 483] is satisfied. It can also be verified that $0 \phi$ satisfies the relevant scalar field equation appearing on [6] p. 483]. Since $\text{tr} \mathcal{K}^2 + \Phi_0^2 = 1$, [6] (10a), p. 483] holds. What remains is to demonstrate that [6] (10b), p. 483] is fulfilled. Note that this equality can be rewritten as
\[ \nabla_A 0 K_B^A = 0 j_B, \]
where $0 K_B^A = 0 k_B^A$ (we use the capital letter $K$ to distinguish the Weingarten map from the second fundamental form) and $\nabla$ is the Levi-Civita connection induced by $0 g$. By arguments similar to the proof of Lemma [59] it can be calculated that
\[ \nabla_A 0 K_B^A = - t^{-1} \mathcal{Y}_{A}(p_B) + \sum A t^{-1}(p_A - p_B) \mathcal{Y}_A(\nabla A \mathcal{Y}_B), \]
\[ \mathcal{Y}_A(\nabla A \mathcal{Y}_B) = \mathcal{Y}_A(\tilde{D}_A \mathcal{Y}_B) + \mathcal{Y}_B(p_A) \ln t, \]
where there is no summation on $B$ in the first equality and no summation on any index in the second equality. Moreover, $\tilde{D}$ here denotes the Levi-Civita connection induced by $\tilde{h}$; note that $\{ \mathcal{Y}_A \}$ is an orthonormal frame with respect to $\tilde{h}$. Due to these observations, it follows that
\[ \nabla_A 0 K_B^A = - t^{-1} \tilde{D}_A \mathcal{K}_B^A + \sum A t^{-1}(p_A - p_B) \mathcal{X}_B(p_A) \ln t. \]
The second term on the right hand side can be calculated to be
\[ \sum A t^{-1}(p_A - p_B) \mathcal{X}_B(p_A) \ln t = \frac{1}{2} t^{-1} \mathcal{Y}_B (\sum A p_A^2) \ln t - t^{-1} p_B \mathcal{X}_B (\sum A p_A) \ln t. \]
However, $\sum A p_A = 1$, so that the second term on the right hand side vanishes. Combining this observation with $\sum A p_A^2 + \Phi_0^2 = 1$ yields
\[ \sum A t^{-1}(p_A - p_B) \mathcal{X}_B(p_A) \ln t = - t^{-1} \ln t \cdot \Phi_a \mathcal{X}_B(\Phi_a). \]
Combining this equality with \((73)\) and the fact that $\text{div}_h \mathcal{K} = \Phi_a d \Phi_a$ yields
\[ \nabla_A 0 K_B^A = - t^{-1} \Phi_a \mathcal{X}_B(\Phi_a) - t^{-1} \Phi_a \mathcal{X}_B(\Phi_a). \]
On the other hand,
\[ 0 j_B = - t^{-1} \ln t \cdot \Phi_a \mathcal{X}_B(\Phi_a) - t^{-1} \Phi_a \mathcal{X}_B(\Phi_a). \]
To conclude, [6] (10b), p. 483] holds. In short, $(0 g, 0 k, 0 \phi)$ is a solution to the velocity dominated Einstein-scalar field equations on $\bar{M} \times (0, \infty)$ in the sense of [6].

Next, we prove Proposition [21].

**Proposition [21]**. That $\text{tr}^0 k$ can be assumed to equal $- t^{-1}$ and that $- t \cdot 0 K$ is independent of time follows from the discussion in [6] Subsection 2.2, pp. 482–484]. Next, since $\mathcal{K}$ is symmetric with respect to $0 g$, the eigenvalues of $\mathcal{K}$ are real and the local frame $\{ \mathcal{X}_A \}$ is orthogonal with respect to $0 g$. Moreover, it can be verified that $t^{-2} p_A^0 g_{AB}$ (no summation) is independent of $t$. We can therefore define $\alpha_A$ to be the strictly positive functions such that $\sum A = 0 g_{AA} |_{t=1}$ (no summation). This means that $0 g$ takes the form \((20)\). Define $\mathcal{X}_A$ by $\mathcal{X}_A := \alpha_A^{-1} \mathcal{X}_A$ and let $\{ \mathcal{Y}_A \}$ be the frame dual to $\{ \mathcal{X}_A \}$. Then $0 g$ and $0 k$ take the form \((19a)\) and \((19b)\) respectively. Finally, since $\partial_t (0 \phi) = 0$, cf. [6] (15), p. 484], $0 \phi$ takes the form \((19c)\).

Next, let us verify that the conditions of Definition [10] are satisfied. That $\text{tr} \mathcal{K} = 1$ follows from the fact that $\text{tr}^0 k = - t^{-1}$. That $\mathcal{K}$ is symmetric with respect to $\tilde{h}$ follows from the fact that both are diagonal with respect to the frame $\{ \mathcal{X}_A \}$. That the eigenvalues of $\mathcal{K}$ are distinct is
an assumption. Since $1 + p_1 - p_{n-1} - p_n > 0$, the final condition of Definition \[10\] is satisfied. Combining \[6\] (10a) and (12a), p. 483, keeping Remark \[19\] in mind, it follows that

$$-t^{-2} \text{tr} X^2 + t^{-2} = t^{-2} \Phi_b^2.$$ 

In particular, $\text{tr} X^2 + \Phi_b^2 = 1$. What remains is to verify that $\text{div}_t X = \Phi_a d\Phi_b$. Note, to this end, that \[73\] and \[74\] are satisfied in the current situation. This means that

$$\nabla_A^0 K^A_B = -t^{-1} D_A X^B_B - t^{-1} \ln t \cdot \Phi_a \cdot \nabla_B(\Phi_a).$$

On the other hand, the equality \[75\] is also satisfied in the current situation. Since $\nabla_A^0 K^A_B = 0^{jB}$, cf. \[6\] (10b), p. 483, it follows that $\text{div}_t X = \Phi_a d\Phi_b$. The proposition follows. \[\square\]

Finally, we prove Theorem \[22\]. Due to Proposition \[18\] $(0^c, \Phi, 0, 0)$ defined by \[19\] is a solution to the velocity dominated Einstein-scalar field equations on $M \times (0, \infty)$ in the sense of \[6\], keeping Remark \[19\] in mind (as in the proof of Proposition \[15\]), we focus on the case $n = 3$ and therefore refer to \[6\]; the general case follows by appealing to \[22\], in particular \[22\] Theorem 10.2, p. 1104 and \[22\] Remark 1, p. 1105). Assuming that the data $(\dot{M}, \dot{h}, \dot{X}, \Phi_a, \Phi_b)$ are real analytic and that the eigenvalues satisfy $1 + p_1 - p_{n-1} - p_n > 0$ (i.e., $p_A > 0$, $A = 1, 2, 3$, in case $n = 3$), we can apply \[6\] Theorem 1, p. 484. This yields a real analytic solution to the Einstein-scalar field equations on a set of the form $M$ introduced in the statement of the theorem. Moreover, on a fixed compact subset of $M$, say $K$, this solution has asymptotics of the form \[6\] 1-6, Theorem 1, p. 484, where the estimates are uniform on $K$. Finally, the solution is unique with these properties.

In order to relate the results of \[6\] with the conclusions we wish to draw here, note that \[6\] Theorem 1, p. 484) yields

$$0^c g^{AC} g_{CB} = \delta_B^A + o(t^{p_A}),$$

where $\alpha^A_B$ are strictly positive numbers and the indices refer to the frame $\{X_A\}$ and co-frame $\{\nabla^A\}$. On the other hand, representing the metric as in \[13\] yields

$$0^c g^{AC} g_{CB} = t^{-2p_A} g_{AB} = t^{2(p_{\text{max}}(A, B) - p_A)b_{AB}}.$$ 

Combining these two equalities, it follows that the diagonal components of $b_{AB}$ converge to 1 at a rate. Moreover, if $A > B$, we conclude that $b_{AB}$ converges to 0 at a rate. Since $b$ is symmetric, the same holds for $b_{BA}$. This means that $\dot{h}$, defined by \[14\], converges to $\dot{h}$ at a rate and uniformly on compact subsets; i.e., \[21\] holds. Combining this observation with arguments similar to the ones presented in connection with \[76\] below yields the conclusion that $\dot{q}$ and $\ln t$ are interchangeable; cf. \[78\] below. Next, note that \[6\] Theorem 1, p. 484) yields the conclusion that $\theta h$ converges to 1 at a rate, and that $\theta - K$ converges to zero at a rate; i.e., \[21\] holds. The conclusions concerning the scalar field, i.e., \[21c\] and \[21d\], follow from \[6\] Theorem 1, p. 484, the fact that $\dot{q}$ and $\ln t$ are interchangeable, the fact that $\theta h$ converges to 1 at a rate and the fact that $\dot{U} = \theta^{-1} \partial_t$. \[\square\]

5. IMPROVING THE ASYMPTOTICS IN THE GENERALISED SETTING

The purpose of the present section is to provide a proof of Theorem \[27\] We proceed gradually, starting with estimates of the relative spatial variation of the mean curvature. We continue with estimates of the eigenvector fields of the expansion normalised Weingarten map as well as of the components of the expansion normalised Weingarten map with respect to the eigenvector fields of $\nabla^\prime$. Given this information, we are in a position to estimate $\theta^{-2} \dot{S}$, where $S$ is the spatial scalar curvature, and $\theta^{-2} \dot{R}$, where $\dot{R}$ is the spatial Ricci tensor, viewed as a $(1,1)$-tensor field. We also obtain estimates for similar quantities involving the lapse function. Combining this information, we are in a position to estimate the deceleration parameter $q$ as well as the expansion normalised normal derivative of $K$, $\dot{L}_Y K$. We also obtain more detailed estimates of the components of $\dot{L}_Y K$ with respect to the eigenvector frame associated with $X$. Given the above control of the geometry, it is possible to appeal to the results of \[65\] in order to control the asymptotics of the
scalar field. In the final step, we then combine all of the above in order to conclude that the equations corresponding to the limits of the Hamiltonian and momentum constraints hold.

5.1. Estimates for the relative spatial variation of the mean curvature. Consider a development which is locally (close to the singularity) either of the form \([13]\) (in general dimensions) or of the form \([17]\). For such a development, define the logarithmic volume density, \(\varrho\), by the relation \(e^\varrho \mu_{\bar{h}} = \mu_{\bar{g}}\), where \(\bar{h}\) is the metric appearing in the definition of the initial data and \(\bar{g}\) is the metric induced on hypersurfaces of the form \(\bar{M}_t := \bar{M} \times \{ t \}\). This means that

\[
e^\varrho = e^\varrho|\mu_{\bar{h}}(\mathcal{X}_1, \ldots, \mathcal{X}_n)| = |\mu_{\bar{g}}(\mathcal{X}_1, \ldots, \mathcal{X}_n)|,\]

assuming \(|\mathcal{X}_A|_{\bar{h}} = 1, A = 1, \ldots, n\). In the case of \([17]\), taking the square of this equality yields

\[
e^{2\varrho} = \theta^{-2} \begin{vmatrix} b_{11} & b_{12} & \cdots & b_{1n} \\ b_{21} & b_{22} & \cdots & b_{2n} \\ \vdots & \vdots & \cdots & \vdots \\ b_{n1} & b_{n2} & \cdots & b_{nn} \end{vmatrix},\]

note that the \(p_i\) should here be thought of as given smooth functions satisfying \(p_1 < \cdots < p_n\). In the case of \([13]\), this equality holds with \(\theta\) replaced by \(t^{-1}\). In what follows, we assume \(\bar{h}\) to converge to \(\bar{h}\) in \(C^i\) for any \(0 \leq l \in \mathbb{N}\) at a rate \(\theta^{-\epsilon}\) in the case of \([17]\) and at a rate \(t^{-\epsilon}\) in the case of \([13]\). This means that \(b_{AB}\) converges to \(\delta_{AB}\) at this rate. In particular, in the case of \([17]\),

\[
\|\theta + \ln \theta\|_{C^i(\bar{M})} \leq C_{\ell} \theta^{-\epsilon}
\]

on \(M_0 := \bar{M} \times (0, t_0]\), where \(t_0 \in (0, t_+).\) Similarly, in the case of \([13]\), there is, for each \(0 \leq l \in \mathbb{N}\), a constant \(C_l\) such that

\[
\|\theta + \ln \theta\|_{C^l(\bar{M})} \leq C_l t^{2(\epsilon - 1)}
\]

for all \(t \leq t_0\). The reason we obtain \(C^l\) bounds in \([78]\) and only \(C^0\) bounds in \([77]\) is that we, of course, control the spatial derivatives of \(t\), but not, a priori, the spatial derivatives of \(\theta\). In order to remedy this situation, we first derive bounds on the spatial derivatives of \(\varrho\), given suitable assumptions concerning \(\bar{N}\). These bounds can then be used to derive estimates for \(\ln \theta\).

Lemma 61. Consider a metric of the form \([22]\) and assume that there, for every \(1 \leq l \in \mathbb{N}\), are constants \(C_l\) and \(a_l\) such that for \(1 \leq k \leq l\),

\[
|\bar{D}^k \ln \bar{N}|_{\bar{h}} \leq C_l (\ln \theta)^{a_l}
\]

on \(M_0\), where \(\bar{N} = \theta N\). Assume, moreover, that \([88]\) holds and that \(\theta\) diverges uniformly to \(\infty\) as \(t \to 0^+\). Then there are, for every \(1 \leq l \in \mathbb{N}\), constants \(K_l\) and \(b_l\) such that

\[
|\bar{D}^k \varrho|_{\bar{h}} \leq K_l (\ln \theta)^{b_l}, \quad |\bar{D}^k \theta|_{\bar{h}} \leq K_l (\theta)^{b_l},
\]

\[
|\bar{D}^k \ln \theta|_{\bar{h}} \leq K_l (\ln \theta)^{b_l}, \quad |\bar{D}^k \ln \theta|_{\bar{h}} \leq K_l (\theta)^{b_l}
\]

on \(M_0\) for \(k \leq l\).

Remark 62. Given \([81]\), one can return to \([76]\) in order to deduce that, for every \(0 \leq l \in \mathbb{N}\), there is a constant \(C_l\) such that

\[
|\bar{D}^k (\theta + \ln \theta)|_{\bar{h}} \leq C_l \theta^{-\epsilon}
\]

on \(M_0\) for all \(k \leq l\).

Remark 63. If the conditions of Theorem \([27]\) are satisfied, then the conditions of the lemma are satisfied.

Proof. Let \(\{E_i\}\) be a local orthonormal frame with respect to \(\bar{h}\). In what follows, we then use the notation \(I, E_1 = E_{I_1} \cdots E_{I_m}\) and \(|I| = m\) for vectors of the form \(I = (I_1, \ldots, I_m)\), where \(I_j \in \{1, \ldots, n\}\).
Due to \[76\] (7.9), p. 74, \( \partial_t \dot{\varrho} = \dot{N} \), so that \( \partial_t E_1 \varrho = E_1 \dot{N} \). Due to this equality, it follows that \( \dot{U}(E_1 \varrho) \) can be written as a linear combination of terms of the form \( E_1, \ln \dot{N}, \cdots E_k \ln \dot{N} \), where \( |I_1| + \cdots + |I_k| = |I| \). Letting \( \gamma \) be an integral curve of \( \dot{U} \), note that
\[
\frac{d}{ds} \varrho \circ \gamma = (\dot{U} \varrho) \circ \gamma = 1;
\]
\text{cf.} \[76\] (7.9), p. 74. We can thus choose the parametrisation of \( \gamma \) in such a way that \( \varrho \circ \gamma(s) = s \). Then
\[
\left| \frac{d}{ds} (E_1 \varrho \circ \gamma)(s) \right| \leq C_I (\ln \varrho \circ \gamma(s))^{c_I} \leq C_I (s)^{c_I},
\]
where we appealed to \[77\], \[79\] and the fact that \( \varrho \circ \gamma(s) = s \). This estimate can be integrated in order to deduce that \[80\] holds. In order to obtain these estimates, we, again, appealed to \[77\] and the fact that \( \varrho \circ \gamma(s) = s \).

In order to estimate the derivatives of \( \ln \theta \), note that \[76\] yields
\[
\exp(2 \varrho + 2 \ln \theta) = b_{11} \cdots b_{nn} + R,
\]
where \( R \), up to signs, consists of terms containing \( k \geq 1 \) factors of the form \( \theta^{-2(p_A - p_B)} b_{AB} \), \( A > B \), and \( n - k \) factors of the form \( b_{AB} \), \( A \leq B \). Differentiating this equality once yields
\[
2 \exp(2 \varrho + 2 \ln \theta) E_i \varrho + 2 \exp(2 \varrho + 2 \ln \theta) E_i \ln \theta = E_i (b_{11} \cdots b_{nn}) + E_i R.
\]
Note that the factors \( \exp(2 \varrho + 2 \ln \theta) \) converge to 1 due to \[77\]. Moreover, the first term on the right hand side converges to zero as \( \theta^{-\varepsilon} \). Consider \( E_i R \). It contains terms with two types of factors. Either \( E_i \) hits \( \theta^{-2(p_A - p_B)} b_{AB} \), \( A > B \), in which case the result is
\[
-2 E_i (p_A - p_B) \ln \theta \cdot \theta^{-2(p_A - p_B)} b_{AB} - 2 (p_A - p_B) E_i (\ln \theta) \theta^{-2(p_A - p_B)} b_{AB} + \theta^{-2(p_A - p_B)} E_i b_{AB}.
\]
The first term converges to zero faster than \( \theta^{-\varepsilon} \), the second term can be written as \( E_i (\ln \theta) \) times a factor that converges to zero faster than \( \theta^{-\varepsilon} \), and the last term converges to zero faster than \( \theta^{-\varepsilon} \). Second, if \( E_i \) hits a factor of the form \( b_{AB} \), \( A \leq B \), then the result converges to zero at least at the rate \( \theta^{-\varepsilon} \). To conclude, the right hand side of the equality \[84\] contains a term including the factor \( E_i \ln \theta \). However, in this term, the factor multiplying \( E_i \ln \theta \) converges to zero faster than \( \theta^{-\varepsilon} \). The term on the right hand side including the factor \( E_i \ln \theta \) can thus be combined with the second term on the left hand side. This yields a bound of \( E_i \ln \theta \) in terms of \( E_i \varrho \) and decaying terms. In order to estimate higher order derivatives, one can proceed similarly, differentiating \[84\] further. By an inductive argument, all but the highest order derivatives appearing on the right hand side can be estimated. The highest order derivative of \( \ln \theta \) appearing on the right hand side will, in analogy with the above, be multiplied by a factor converging to zero. This can then be moved to the left hand side, as in the first step. At the end, we then obtain the estimate \[81\]. \( \square \)

5.2. Components of the eigenvector fields. Next, introduce the matrices \( X_A^B \) and \( Y_A^B \) by the requirement that \( X_A = X_A^B \mathcal{X}_B \) and \( \mathcal{X}_A = Y_A^B X_B \), where \( \{ \mathcal{X}_A \} \) and \( \{ X_A \} \) are defined in the statement of Theorem \[27\]. Then
\[
X_A^B = \mathcal{Y}_B(X_A), \quad Y_A^B = Y_B(\mathcal{X}_A), \quad X_A^B Y_A^C = Y_B^A X_C^D = \delta_A^C.
\]
These matrices satisfy the following relations.

Lemma 64. Given the assumptions and notation of Theorem \[27\], there is an \( \eta > 0 \) and, for all \( k \in \mathbb{N} \), a constant \( C_k \), such that
\[
|\tilde{D}^k (X_A^B - \delta_A^B)|_k \leq C_k \theta^{-2n} \min \{ 1, \theta^{-2(p_B - p_A)} \},
\]
\[
|\tilde{D}^k (Y_A^B - \delta_A^B)|_k \leq C_k \theta^{-2n} \min \{ 1, \theta^{-2(p_B - p_A)} \},
\]
\[
|\tilde{D}^k (\tilde{\mu}_A + p_A \ln \theta)|_k \leq C_k \theta^{-2n}
\]
on \( M_0 \) for all \( A \) and \( B \).
Proof. Note, to begin with, that \((89)\) and the definition of \(X_A, Y^A, \mathcal{Y}_A\) and \(\mathcal{Y}^A\) yield the conclusion that for all \(k \in \mathbb{N}\), there are constants \(C_k\) and \(\eta > 0\) such that
\[
|\tilde{D}^k[\mathcal{Y}^A(X_B) - \delta^A_B]|_h + |\tilde{D}^k[Y^A(\mathcal{Y}_B) - \delta^A_B]|_h + |\tilde{D}^k(b_{AB} - \delta_{AB})|_h \leq C_k\theta^{-2\eta}
\]
on \(M_0\) for all \(A, B\). In particular, \((86)\) and \((87)\) hold in case \(B \leq A\). On the other hand, \((90)\)
\[
\sum_{A, B} b_{AB} \theta^{-2p_{\max}(A, B)} [\mathcal{Y}^A \otimes \mathcal{Y}^B] = \sum_{C} e^{2\tilde{\mu} C} Y_C \otimes Y_C.
\]
Evaluating this equality on \((\mathcal{Y}_A, X_A)\) yields
\[
\sum_{B} b_{AB} \theta^{-2p_{\max}(A, B)} [\mathcal{Y}^B(X_A)] = e^{2\tilde{\mu} A} Y_A(\mathcal{Y}_A)
\]
(no summation on \(A\)). Combining this equality with \((81)\) and \((89)\) yields \((88)\). Next, evaluating \((90)\) on \((\mathcal{Y}_A, X_C)\) with \(C < A\) yields
\[
\sum_{B} b_{AB} \theta^{-2p_{\max}(A, B)} [\mathcal{Y}^B(X_C)] = e^{2\tilde{\mu} C} Y_C(\mathcal{Y}_A).
\]
Combining this equality with \((88)\) and \((89)\) yields \((87)\).

In order to prove \((86)\), note that \(YX = \text{Id}\), where \(X\) and \(Y\) are the matrices with components \(X_A^B\) and \(Y_B^A\) respectively. Next, note that \(Y = U + L\), where \(U\) is upper triangular, \(L\) is lower triangular and the diagonal components of \(L\) vanish. Moreover, for \(A > B\), \(L_B^A = O(\theta^{-2p_A + p_B + \eta})\) in any \(C^m\)-norm. Moreover, \(U - \text{Id} = O(\theta^{-2\eta})\) in any \(C^m\)-norm. In particular,
\[
X = Y^{-1} = (\text{Id} + U^{-1}L)^{-1}U^{-1};
\]
note that \(\|U^{-1}L\| \to 0\) in the direction of the singularity, so that we can assume \(\|U^{-1}L\| \leq 1/2\) in the arguments to follow. Let \(R := -U^{-1}L\). Then
\[
X = \sum_{k=0}^{\infty} R^k U^{-1}.
\]
Note that \(R^k = O(\theta^{-2k\eta})\) in any \(C^m\) norm. Moreover, for \(k\) large enough, say \(k \geq K\), \(2k\eta \geq 2(p_n - p_1 + \eta)\) on \(M\). It is therefore straightforward to prove the desired estimate for the terms with \(k \geq K\), and we only need to focus on a finite number of terms of the form \(R^kU^{-1}\). However, it can be verified that the components of \(R\) satisfy \(R_A^B = 0\), \(A = 1, \ldots, n\), and, for \(B < n\),
\[
R_A^B = \begin{cases} O(\theta^{-2(p_B + p_A + \eta)}) & B \geq A, \\ O(\theta^{-2(p_B - p_A - \eta)}) & B < A. \end{cases}
\]
Moreover, this holds in any \(C^m\) norm. It can be verified that this structure is preserved when taking positive integer powers. Finally, when multiplying a matrix with this structure on the right with \(U^{-1}\) (which is upper triangular), the lower triangular part of the resulting matrix satisfies the same conditions as the lower triangular part of \(R\) (i.e., the estimates satisfied by \(R_A^B\), \(B < A\), are also satisfied by \((R^mU^{-1})_B^A\) for \(m \geq 1\) and \(B < A\)). It follows that \((86)\) holds for \(B > A\). The lemma follows. \(\square\)

Using this information, we are in a position to prove that \((26b)\) holds.

Lemma 65. Given the assumptions and notation of Theorem \(27\), there is an \(\eta > 0\) and, for all \(k \in \mathbb{N}\), a constant \(C_k\), such that
\[
|D^k[K(\mathcal{Y}^A, \mathcal{Y}_B) - p_B^A]|_h \leq C_k e^{-2\eta} \min\{1, \theta^{-2(p_B - p_A)}\}
\]
(no summation on \(B\)) on \(M_0\).

Proof. Note, to begin with, that
\[
K_B^A := K(\mathcal{Y}^A, \mathcal{Y}_B) = \sum_C \ell_C \mathcal{Y}^A(X_C) Y_C(\mathcal{Y}_B) = \sum_C \ell_C X^A_C Y_C^B,
\]
where \(\ell_A\) is the eigenvalue of \(K\) corresponding to \(X_A\). Note also that \(\ell_A - p_A = O(\theta^{-2\eta})\) in any \(C^m\)-norm due to \((18a)\). Combining these observations with Lemma \(64\) yields \(K_B^A = p_A + O(\theta^{-2\eta})\) (no summation on \(A\)) in any \(C^k\)-norm. Similarly, \(K_B^A = O(\theta^{-2\eta})\) for \(B < A\) in any \(C^k\)-norm. Finally, for \(A < B\), the desired conclusion again follows by appealing to Lemma \(64\). The desired estimate follows. \(\square\)
Next, define $\lambda_{BC}^A$ by $[X_B, X_C] = \lambda_{BC}^A X_A$. We wish to estimate $\lambda_{BC}^A$ in case $1 + p_A - p_B - p_C \leq 0$; cf. the definition of quiescent initial data on the singularity.

**Lemma 66.** Given the assumptions and notation of Theorem 27, fix $\bar{x} \in \bar{M}$ and assume that $1 + p_A(\bar{x}) - p_B(\bar{x}) - p_C(\bar{x}) \leq 0$ for some $A$ and $B \neq C$. Then there is a neighbourhood $U$ of $\bar{x}$ in $\bar{M}$, an $\eta > 0$ and, for every $k \in \mathbb{N}$, a constant $C_k$ such that

$$|D^k(\theta^{-2(1+p_A-p_B-p_C)}\lambda_{BC}^A)|_h \leq C_k \theta^{-2\eta}$$

on $U \times (0, t_0)$.

**Remark 67.** Assuming, in addition to the conditions of the lemma, that $n = 3$ and $\mathcal{E}_{23}^1 = 0$ on $\bar{M}$, there is an $\eta > 0$ and, for all $k \in \mathbb{N}$, a constant $C_k$, such that

$$|D^k(\lambda_{23}^1)|_h \leq C_k \theta^{-2(p_2-p_1+\eta)}$$

on $M_0$.

**Proof.** Compute

$$Y^A([X_B, X_C]) = X_B^D \mathcal{F}_D(X_C^E)Y_E^A - X_C^D \mathcal{F}_D(X_B^E)Y_E^A + X_B^D X_C^F \gamma_{DE},$$

where $\gamma_{DE} = \mathcal{F}_D([\mathcal{F}_D, \mathcal{F}_E])$. Assume that $1 + p_A(\bar{x}) - p_B(\bar{x}) - p_C(\bar{x}) \leq 0$. In order for this to hold, we have to have $B > A$ and $C > A$; cf. Remark 30. In what follows, we also assume, without loss of generality, that $B < C$. Combining these observations with Lemma 63, it is clear that the first two terms on the right hand side of (94) are $O(\theta^{-2(p_B-p_A+\eta)})$ in any $C^k$ norm. In particular, multiplying the first two terms on the right hand side of (94) with $\theta^{-2(1+p_A-p_B-p_C)}$, it is clear that the resulting expression is $O(\theta^{-2(1-p_B+\eta)})$ in any $C^k$-norm. Next, consider

$$\theta^{-2(1+p_A-p_B-p_C)}X_B^D X_C^F \gamma_{DE} = (\theta^2(p_B-p_D)X_B^D) \cdot (\theta^2(p_C-p_E)X_C^F) \cdot (\theta^2(p_F-p_A)Y_A^E) \cdot \theta^{-2(1+p_B-p_D-p_E)} \gamma_{DE}.$$

Due to Lemma 64, the first three factors on the right hand side are bounded in any $C^k$ norm. In order for the last term to be non-zero, $1 + p_F - p_D - p_E$ has to be strictly positive in a neighbourhood of $\bar{x}$. The estimate (92) follows. The proof of (93) is similar. \hfill \Box

**Corollary 68.** Given the assumptions and notation of Theorem 27, there is an $\eta > 0$ and, for every $k \in \mathbb{N}$, a constant $C_k$ such that

$$|D^k(\theta^{-2(1+p_A-p_B-p_C)}\lambda_{BC}^A)|_h \leq C_k \theta^{-2\eta}$$

on $M_0$ for all $A, B$ and $C$.

**Remark 69.** Since $\lambda_{BC}^A$ is well defined up to a sign, the estimate (95) is globally meaningful.

**Proof.** If $B = C$, there is nothing to prove. In what follows, we therefore assume $B \neq C$. Fix an $\bar{x} \in \bar{M}$. Given $A, B, C$ such that $B \neq C$, there are then two possibilities. Either $1 + p_A(\bar{x}) - p_B(\bar{x}) - p_C(\bar{x}) \leq 0$ or $1 + p_A(\bar{x}) - p_B(\bar{x}) - p_C(\bar{x}) > 0$. In the first case, we can appeal to Lemma 66 in order to conclude that there is an open neighbourhood $U_{ABC}$ and an $\eta_{ABC} > 0$ such that (92), with $\eta$ replaced by $\eta_{ABC}$, holds in $U_{ABC} \times (0, t_0]$. In case $1 + p_A(\bar{x}) - p_B(\bar{x}) - p_C(\bar{x}) > 0$, there is trivially such an open neighbourhood $U_{ABC}$ and $\eta_{ABC} > 0$. Taking the intersection of all these $U_{ABC}$ yields an open neighbourhood $U_{\bar{x}}$ and taking the minimum of all the $\eta_{ABC}$ yields an $\eta_{\bar{x}} > 0$. Since $\bar{M}$ is compact, there is a finite cover by sets of the form $U_{\bar{x}}$. Taking the minimum of the associated $\eta_{\bar{x}}$ yields an $\eta > 0$ such that the statement of the corollary holds. \hfill \Box

Using this information, we estimate the conformally rescaled spatial Ricci curvature.

**Proposition 70.** Assume that the conditions of Theorem 27 are satisfied. Let, using the notation of Theorem 27, $\mathcal{R}$ denote the spatial Ricci tensor, viewed as a family of $(1,1)$-tensor fields, $\mathcal{S}$ denote the spatial scalar curvature, $\mathcal{R} := \theta^{-2} \mathcal{R}$ and $\mathcal{S} := \theta^{-2} \mathcal{S}$. Then there is an $\eta > 0$ and, for each $k \in \mathbb{N}$, a $C_k$ such that

$$|\mathcal{D}^k \mathcal{S}|_h + |\mathcal{D}^k \mathcal{R}|_h \leq C_k \theta^{-2\eta}$$
Remark 71. If, in addition to the assumptions of the proposition, \( n = 3 \) and \( \mathcal{O}_{23}^1 = 0 \) on \( M \), then there are, for each \( k \in \mathbb{N} \), constants \( a_k \) and \( C_k \) such that
\[
|\bar{D}^k \bar{S}|_h + |\bar{D}^k \bar{R}|_h \leq C_k (\ln \theta)^{a_k} \theta^{-2(1 - p_A)}
\]
on \( M_0 \).

Proof. To estimate \( \bar{R} \), it is convenient to appeal to [66, Corollary 96, p. 63]. The \( \mu_A \) appearing in this corollary are given by \( \mu_A = \bar{\mu}_A + \ln \theta \), so that \( \bar{D}^k [\mu_A - (1 - p_A) \ln \theta] = O(\theta^{-2\eta}) \) due to (88).

Note that \( \lambda^A_{BC} \) (the \( \gamma^A_{BC} \) appearing in [66, Corollary 96, p. 63] should be replaced by \( \lambda^A_{BC} \)) here) and \( a_A \) are bounded in any \( C^k \)-norm and that \( \bar{a}_A, \bar{\mu}_A \) and \( \mu_A \) do not grow faster than polynomially in \( \ln \theta \) in any \( C^k \)-norm. In particular, any number of \( X_A \) derivatives of these objects grows at worst polynomially. Note also that \( p_A < 1 \) due to Remark [50]. Combining these observations with [66, (221)–(226), p. 63], it is clear that for \( k \in \mathbb{N} \), there is a constant \( \alpha_k \) such that most of the terms on the right hand side of [66, (221), p. 63] decay to zero as \( O((\ln \theta)^{\alpha_k} \theta^{-2(1 - p_A)}) \) in \( C^k \). However, there are some terms containing factors of the form \( e^{2\mu_A - 2\mu_B - 2\mu_C} \). On the other hand, terms that contain a factor \( e^{2\mu_A - 2\mu_B - 2\mu_C} \) also contain a factor \( \lambda^A_{BC} \) (or a spatial derivative thereof). In order to estimate the corresponding terms, we can thus appeal to Corollary [65].

In order to illustrate the arguments required, consider, e.g.,
\[
e^{2\mu_A - 2\mu_B - 2\mu_C} X_D (\lambda^A_{BC})
= e^{2\mu_A - 2\mu_B - 2\mu_C} \theta^{2(1 + p_A - p_B - p_C)} X_D (\theta^{-2 - 2(1 + p_A - p_B - p_C)} \lambda^A_{BC})
+ e^{2\mu_A - 2\mu_B - 2\mu_C} \theta^{2(1 + p_A - p_B - p_C)} X_D [2(1 + p_A - p_B - p_C) \ln \theta] (\theta^{-2 - 2(1 + p_A - p_B - p_C)} \lambda^A_{BC}).
\]

Due to Corollary [66] terms of this form decay as desired in any \( C^k \)-norm; since \( \mu_A - (1 - p_A) \ln \theta = O(\theta^{-2\eta}) \), the product of the first two factors is, in each term, bounded in any \( C^k \)-norm; the third factor in the second term does not grow faster than polynomially in any \( C^k \)-norm; and the last factor decays as \( \theta^{-2\eta} \) in any \( C^k \)-norm in each of the terms.

In case \( n = 3 \), a more detailed argument can be carried out. In fact, then, if \( 1 \in \{ B, C \} \), \( e^{2\mu_A - 2\mu_B - 2\mu_C} \) decays to zero as \( O((\ln \theta)^{\alpha_k} \theta^{-2(1 - p_A)}) \) in \( C^k \). Assume, therefore, that \( \{ B, C \} = \{ 2, 3 \} \). If \( A \in \{ 2, 3 \} \), we again obtain the same decay. The only interesting terms are thus the ones containing one factor of the form \( e^{2\mu_A - 2\mu_B - 2\mu_C} \) and a factor of the form \( \lambda^1_{23} \) or a spatial derivative thereof. However, we know that \( \lambda^1_{23} = O(\theta^{-2(p_A - p_B - p_C)}) \) (and similarly for every \( C^k \)-norm); cf. Remark [57]. Thus the potentially problematic terms are \( O(\theta^{-2(1 - p_A)}) \) (and similarly for every \( C^k \)-norm).

Writing
\[
\bar{R} = \bar{R}^A_X Y_A \otimes Y_B,
\]
the desired estimate for \( \bar{R} \) follows. Due to [66, Lemma 93, p. 58–59], the proof of the estimate for \( S \) is essentially identical.

Next, we wish to estimate the contribution from the lapse function to \( \tilde{L}_U^K \). Introduce, to this end, the \((1, 1)\)-tensor field \( \mathcal{N} \), given by
\[
\mathcal{N}^A_B := \theta^{-2} N^{-1} \nabla^A \nabla_B N.
\]
Here \( \nabla \) is the Levi-Civita connection induced on the hypersurfaces of constant \( t \).

Lemma 72. Assume that the conditions of Theorem [27] are satisfied and let \( \mathcal{N} \) be defined by (98). Then, using the notation of Theorem [27], there is an \( \eta > 0 \) and, for each \( k \in \mathbb{N} \), a \( C_k \) such that
\[
|\bar{D}^k \mathcal{N}|_h \leq C_k \theta^{-2\eta}
\]
on \( M_0 \).
Remark 73. If, in addition to the assumptions of the proposition, \( n = 3 \) and \( \Theta^1_{23} = 0 \) on \( \bar{M} \), then there are, for each \( k \in \mathbb{N} \), constants \( a_k \) and \( C_k \) such that

\[
|\bar{D}^k \bar{N}|_h \leq C_k (\ln \theta)^{a_k} \theta^{2(1-p_3)}
\]
on \( M_0 \).

Proof. Note, first of all, that the derivatives of \( \ln N \) can be bounded by appealing to \([65, (3.4), p. 26]\). Combining this observation with \([65, \text{ Lemma 98, p. 63–64}]\) and the assumptions, the proof of the statement is similar to, but simpler than, the proof of Proposition \( 70 \). □

Next, we estimate \( \bar{L}_U K \).

Proposition 74. Assume that the conditions of Theorem \( 27 \) are satisfied. Then, using the notation of Theorem \( 27 \), there is an \( \eta > 0 \) and, for each \( k \in \mathbb{N} \), a \( C_k \) such that

\[
|\bar{D}^k[q - (n - 1)|_h + |\bar{D}^k \bar{L}_U K|_h \leq C_k \theta^{-2\eta}
\]
on \( M_0 \).

Remark 75. If, in addition to the assumptions of the proposition, \( n = 3 \) and \( \Theta^1_{23} = 0 \) on \( \bar{M} \), then there are, for each \( k \in \mathbb{N} \), constants \( a_k \) and \( C_k \) such that

\[
|\bar{D}^k[q - 2]|_h + |\bar{D}^k \bar{L}_U K|_h \leq C_k (\ln \theta)^{a_k} \theta^{-2(1-p_3)}
\]
on \( M_0 \).

Remark 76. When appealing to the results of \([65]\), we need upper bounds on the eigenvalues of \( \bar{K} \), the Weingarten map of the conformally rescaled metric \( \bar{g} := \theta^2 g \). Combining \([65, (3.3), p. 26]\) with \([65, (3.4), p. 26]\), we conclude that \( \bar{K} = K - (q + 1) \text{Id}/n \). Since \( q \) converges to \( n - 1 \), cf. \([100]\), we conclude that the eigenvalues of \( \bar{K} \) converge to \( p_A - 1 \). In particular, they are all negative, cf. Remark \( 30 \) so that \( \bar{K} \) has a silent upper bound asymptotically; cf. \([65, \text{ Definition 3.10, p. 27}]\).

Proof. The basis for the estimate of \( \bar{L}_U K \) is the formula for \( \bar{L}_U K \) given by \([66, (65), p. 28]\):

\[
(\bar{L}_U K)^i_j = -\left( \frac{\Delta g N}{\theta^2 N} + \frac{n}{n - 1} \frac{\rho - \bar{p} + 2\Lambda}{\theta^2} - \frac{\bar{S}}{\theta^2} \right) K^i_j + \frac{1}{(n - 1)\theta^2} (\rho - \bar{p} + 2\Lambda) \delta^i_j + \frac{1}{N\theta^2} \bar{S}^i_j \nabla_j N + \theta^{-2}p^i_j - \frac{\bar{R}^i_j}{\theta^2}.
\]

Several of the terms can be estimated by appealing to \([66, 29]\) and the convergence of \( K \). However, there are three types of exceptions: terms including a factor of the form \( \theta^{-2}(\rho - \bar{p}) \), a factor of the form \( \Delta / \theta^2 \) or a factor of the form \( \theta^{-2}P \). Estimates for terms including a factor of the form \( \Delta / \theta^2 \) follow immediately from \([81]\). To estimate the remaining terms, it is sufficient to appeal to \([66, \text{ Remark 84, p. 48}]\); given the corresponding formulae and \([24]\), it is clear that we can proceed as in the previous proofs in order to derive the desired estimate for \( \bar{L}_U K \).

The estimate for \( q - (n - 1) \) can be derived by a similar argument due to \([66, (61), p. 27]\). □

In order to refine this conclusion, it is of interest to consider the components of \( \bar{L}_U K \) with respect to the frame \( \{ \mathcal{F}_\lambda \} \).

Lemma 77. Assume that the conditions of Theorem \( 27 \) are satisfied. Then, using the notation of Theorem \( 27 \), there is a \( \eta > 0 \) and, for each \( k \in \mathbb{N} \), a constant \( C_k \) such that

\[
|\bar{D}^k[(\bar{L}_U K)(\mathcal{F}^A, \mathcal{F}^B)]|_h \leq C_k \theta^{-2\eta} \min\{1, \theta^{-2(p_B - p_A)}\}
\]
on \( M_0 \).

Remark 78. If, in addition to the assumptions of the lemma, \( n = 3 \) and \( \Theta^1_{23} = 0 \) on \( \bar{M} \), then there are, for each \( k \in \mathbb{N} \), constants \( a_k \) and \( C_k \) such that

\[
|\bar{D}^k[(\bar{L}_U K)(\mathcal{F}^A, \mathcal{F}^B)]|_h \leq C_k (\ln \theta)^{a_k} \theta^{-2(1-p_3)} \min\{1, \theta^{-2(p_B - p_A)}\}
\]
on \( M_0 \).
Proof. For $A \geq B$, the estimates are immediate consequences of \[(100)\] and \[(101)\]. We therefore need to focus on the case $B > A$. Define, for $A \neq B$,
\[(104)\]
\[
\mathcal{W}_B^A := \frac{1}{\ell_B - \ell_A} (\hat{\mathcal{L}}_U \mathcal{K})(Y^A, X_B);
\]
cf. \[65\] (6.8), p. 68]. Due to Proposition \[74\] we know that the $C^k$-norm of $\mathcal{W}_B^A$ is bounded by the right hand side of \[(100)\]. On the other hand, \[65\] (6.19), p. 70] implies that $e^{\mu_\Lambda - \bar{\mu}_B} \mathcal{W}_B^A$ is antisymmetric for $A \neq B$. In particular, if $A < B$, then
\[(105)\]
\[
\mathcal{W}_B^A = -e^{2(\bar{\mu}_B - \mu_\Lambda)} \mathcal{W}_A^B.
\]
There is thus an $\eta > 0$ and, for each $k \in \mathbb{N}$, a constant $C_k$ such that
\[(106)\]
\[
|\hat{D}^k \mathcal{W}_B^A|_h + |\hat{D}^k [(\hat{\mathcal{L}}_U \mathcal{K})(Y^A, X_B)]|_h \leq C_k \theta^{-2(p_B - p_A)} \theta^{-2\eta}
\]
on $M_0$ if $A < B$, where we appealed to \[(88)\] and to \[(104)\] in order to estimate the second term on the left hand side of \[(106)\]. In case $n = 3$ and $\vartheta_{23} = 0$ on $M$, this estimate can be improved to
\[(107)\]
\[
|\hat{D}^k \mathcal{W}_B^A|_h + |\hat{D}^k [(\hat{\mathcal{L}}_U \mathcal{K})(Y^A, X_B)]|_h \leq C_k (\ln \theta)^{\alpha_k} \theta^{-2(p_B - p_A)} \theta^{-2(1 - p_3)}
\]
on $M_0$ if $A < B$, where $\alpha_k$ and $C_k$ are constants.

Next, note that
\[
(\hat{\mathcal{L}}_U \mathcal{K})(Y^A, X_B) = X_B^A Y_D^B (\hat{\mathcal{L}}_U \mathcal{K})(Y^C, X_D).
\]
Combining this observation with Lemma \[64\] \[(106)\] and \[(107)\] yields the conclusion of the lemma. \[\square\]

5.3. Scalar field. Next, we wish to analyse the asymptotic behaviour of the scalar field. We do so by appealing to the results of \[65\].

Lemma 79. Assume that the conditions of Theorem \[27\] are satisfied. Then, using the notation of Theorem \[27\], there are $\Phi_a, \Phi_b \in C^\infty(M)$, an $\eta > 0$ and, for each $k \in \mathbb{N}$, a constant $C_k$ such that
\[(108)\]
\[
|\hat{D}^k (\hat{U} \phi - \Phi_a)|_h \leq C_k \theta^{-2\eta},
\]
\[(109)\]
\[
|\hat{D}^k (\phi + \Phi_a \ln \theta - \Phi_b)|_h \leq C_k \theta^{-2\eta}
\]
on $M_0$.

Remark 80. Due to \[(82)\], \[(109)\] can be written
\[(110)\]
\[
|\hat{D}^k (\phi - \Phi_a \theta - \Phi_b)|_h \leq C_k \theta^{-2\eta}
\]
on $M_0$.

Proof. The statement essentially follows by appealing to the results of \[65\]. For this reason, the proof mainly consists of a verification that the corresponding conditions are satisfied. Let us begin by verifying that the basic assumptions, cf. \[65\] Definition 3.27, p. 34] are satisfied.

The basic assumptions. It is clear that by restricting $(M, g)$ close enough to the big bang, it is time oriented, it has an expanding partial pointed foliation and $\mathcal{K}$ is non-degenerate. In what follows, it is, however, important to keep in mind that in \[65\], $\vartheta$, $\bar{\vartheta}$ etc. are defined with respect to the reference metric $\bar{g}_{\text{ref}}$ introduced in \[65\] (and $g_{\text{ref}}$ is typically different from $h$). We therefore here write $\bar{g}_{\text{ref}}$ to indicate the logarithmic volume density defined with respect to $\bar{g}_{\text{ref}}$. On the other hand, the difference $\vartheta - \bar{g}_{\text{ref}}$ is independent of $t$. Due to \[65\] Lemma A.1, p. 201], we can assume $\mathcal{K}$ to have a global frame; this is achieved by considering a finite covering space, if necessary, an operation which does not affect the local conclusions concerning the asymptotics. That $\bar{K}$ has a silent upper bound follows from Remark \[76\]. Due to \[(82)\] and Proposition \[74\] it is clear that $\hat{\mathcal{L}}_U \mathcal{K}$ satisfies a weak off diagonal exponential bound; cf. \[65\] Definition 3.19, p. 29] (in this definition, the requirement that the constants $C_{\mathcal{K}, \text{od}}, G_{\mathcal{K}, \text{od}}$ and $\bar{M}_{\mathcal{K}, \text{od}}$ be strictly positive is not necessary; it is sufficient if they are $\geq 0$). Since $\mathcal{K}$ converges exponentially in any $C^k$-norm and since $\chi = 0$ in our setting, it is clear that \[65\] (3.28), p. 34] and \[65\] (3.29), p. 34] hold. Moreover, by assumption, \[65\] (3.18), p. 32] holds. Thus the basic assumptions are satisfied; cf. \[65\] Definition 3.27, p. 34].
Higher order Sobolev assumptions. Next, let us verify that the higher order Sobolev assumptions are satisfied; cf. [65] Definition 3.28, p. 34. Due to the assumptions concerning the derivatives of ln N and the fact that (\varrho_{\text{ref}}) and (ln \theta) are equivalent, cf. [82], it is clear that ln N satisfies the required bound (in what follows, the \( \hat{u} \) appearing in the definition of the higher order Sobolev assumptions will be allowed to depend on the number of derivatives, say \( l \)). Since \( \hat{U} \ln N \) does not grow faster than polynomially in any \( C^k \)-norm by assumption, and the same is true of \( \hat{U}(\ln \theta) = -(1 + q)/n \) (due to (100)), it is clear that \( \hat{U} \ln N \) does not grow faster than polynomially in any \( C^k \)-norm. Combining these observations with the fact that \( \chi = 0 \) and the fact that (18a), (81) and (100) hold, it follows that all the conditions of [65] Definition 3.28, p. 34 are satisfied. In other words, for any choice of \( l \), there is a \( 0 \leq \nu_i \in \mathbb{R} \) such that the \( (\nu_i, l) \)-Sobolev assumptions are satisfied.

Higher order \( C^k \)-assumptions. Due to the arguments of the previous paragraph, there is, for any choice of \( l \), a \( 0 \leq \nu_i \in \mathbb{R} \) such that the \( (\nu_i, l) \)-supremum assumptions are satisfied; cf. [65] Definition 3.31, p. 35.

The standard assumptions. Since \( \chi = 0 \), it is clear that the conditions of [65] Lemma 3.33, p. 36] are satisfied. In particular, the standard assumptions are thus satisfied; cf. [65] Definition 3.36, p. 36.

Basic energy estimate. Next, we wish to appeal to [65] Proposition 14.24, p. 158]. Note, to this end, that, due to the above observations and the assumptions, the conditions of [65] Lemma 7.5, p. 76], [65] Lemma 7.12, p. 81] and [65] Lemma 7.13, p. 83] are all satisfied. Since we are here interested in the wave equation, the estimates [65] (3.32), p. 35] and [65] (3.34), p. 35] are trivially satisfied. Next, note that if \( \kappa_1 \) is the smallest integer strictly larger than \( n/2 + 1 \), then, for a suitable choice of \( u \), the \( (u, \kappa_1) \)-supremum assumptions are satisfied. Given \( 1 \leq l \in \mathbb{Z} \) and a suitable choice of \( u \), the \( (u, l) \)-Sobolev assumptions are also satisfied. Since \( \phi \) satisfies the wave equation, it satisfies [65] (12.32), p. 124] with vanishing right hand side, \( \hat{\chi}^\theta = 0 \), \( \hat{\chi}^\Lambda = 0 \) and \( \hat{\alpha} = 0 \). Due to Proposition 74, it is also clear that [65] (7.78), p. 84] holds. Since \( \alpha \) (using the notation of [65] (1.1), p. 7]) vanishes, it is clear that [65] (11.7), p. 110] holds. That [65] (11.39), p. 115] holds is trivial due to the fact that the left hand side of this estimate vanishes identically. Summing up, the conditions of [65] Proposition 14.24, p. 158] are satisfied. In particular, for any \( 1 \leq l \in \mathbb{Z} \), there are constants \( a_i \) and \( C_i \) such that

\[
\hat{G}_l(\tau) \leq C_l(\tau)^{a_l}
\]

for \( \tau \leq 0 \). Here \( \tau(t) := \varrho_{\text{ref}}(\bar{x}_0, t) \) for some reference point \( \bar{x}_0 \in \bar{M} \). Moreover,

\[
\mathcal{E}_k := \frac{1}{2} \sum_{|I| \leq k} \left( |\hat{U}(E_I \phi)|^2 + \sum_A e^{-2 \mu_A} |X_A(E_I \phi)|^2 + \langle \tau \rangle^{-3} |E_I \phi|^2 \right)
\]

\[
\hat{G}_k(\tau) := \int_{\bar{M}_r} \mathcal{E}_k \mu_{\bar{g}_{\text{ref}}}.
\]

Here \( \{E_i\} \) is a global frame on the tangent space of \( \bar{M} \) and if \( \bar{I} = (i_1, \ldots, i_k) \) is such that \( i_j \in \{1, \ldots, n\} \), then \( E_{\bar{I}} := E_{i_1} \cdots E_{i_k} \) and \( |\bar{I}| = k \). One immediate consequence of (111) is that \( \phi \) does not grow faster than polynomially in any \( C^k \)-norm. However, this is already part of the assumptions. More interestingly, combining this estimate with [69] Corollary 13.8, p. 130] and [65] Remark 13.9, p. 130] yields the conclusion that \( \hat{U} \phi \) does not grow faster than polynomially in any \( C^k \)-norm. Combining this estimate with [66] (167), p. 50] and arguments similar to those of previous lemmas, it follows that \( \hat{U}^2 \phi \) decays as \( \theta^{-2\eta} \) in any \( C^k \)-norm.

Next, note that

\[
\partial_\tau [\hat{U}(\phi)] = \xi \hat{U}^2(\phi),
\]

where \( \xi := \hat{N}/\partial_\tau \tau \). Since \( \xi \) is bounded due to [65] (7.86), p. 85], and since the spatial derivatives of \( \ln \hat{N} \) do not grow faster than polynomially, it is clear that \( \xi \) does not grow faster than polynomially in any \( C^k \)-norm. Thus \( \partial_\tau [\hat{U}(\phi)] \) decays exponentially in any \( C^k \)-norm. In particular, \( \hat{U}(\phi) \)
converges exponentially in any $C^k$-norm. This guarantees the existence of a $\Phi_a \in C^\infty(\tilde{M})$ such that (108) holds. Next, consider
\[(114) \quad \hat{U}(\phi - \Phi_a \theta) = \hat{U}(\phi) - \Phi_a,\]
where we appealed to (65) (7.9), p. 74. Since the right hand side converges to zero exponentially in every $C^k$-norm, we conclude the existence of a $\Phi_b \in C^\infty(\tilde{M})$ such that (110), and thereby (109), holds.

5.4. The limits of the constraint equations in the case of a scalar field. Next, we calculate the limits of the constraint equations in the presence of a scalar field.

**Lemma 81.** Assume that the conditions of Theorem 27 are satisfied. Then, using the notation of Theorem 27, the limits of the renormalised Hamiltonian and momentum constraints can be written
\[(115) \quad \Phi_a^2 + \sum A p_A^2 = 1, \quad \text{div}_b \mathcal{K}' = \Phi_a d\Phi_b\]
respectively, where $\Phi_a$ and $\Phi_b$ are the functions whose existence is guaranteed by Lemma 79.

**Proof.** Due to (66) (7), p. 6, the Hamiltonian constraint implies that
\[(116) \quad 1 = 2\Omega + 2\Omega_\Lambda + \text{tr}K^2 - \theta^{-2}\tilde{S},\]
where $\Omega = \rho/\theta^2$, $\rho$ denotes the energy density associated with the scalar field, $\Omega_\Lambda = \Lambda/\theta^2$ and $\tilde{S}$ denotes the spatial scalar curvature. It is clear that $\Omega_\Lambda$ converges to zero. Moreover, due to (53), p. 48 and the conclusions of Lemma 79 it is clear that $2\Omega$ converges to $\Phi_a^2$. Due to (96), it is also clear that $\theta^{-2}\tilde{S}$ converges to zero. Finally, $\text{tr}K^2$ equals the sum of the squares of the $\ell_A$. This converges to the sum of the squares of the $p_A$. To conclude, the limit of the expansion normalised Hamiltonian constraint is the first equality in (115).

Next, note that the momentum constraint can be written
\[\nabla_m \kappa^m_j + \kappa^m_j \nabla_m \ln \theta - \nabla_j \ln \theta = \hat{U}(\phi)\nabla_j \phi.\]
Compute
\[\nabla_A \kappa^B_C = X_A (\kappa^B_C) - \kappa(\nabla_X A Y^B, X_C) - \kappa(Y^B, \nabla_X A X_C).\]
Choosing $B = A$ and summing over $A$, the first term on the right hand side becomes $X_C (\ell_C)$ (no summation). Next,
\[-\kappa(\nabla_X A Y^A, X_C) = -\ell_C (\nabla_X A Y^A)(X_C) = \ell_C Y^A (\nabla_X A X_C)\]
(no summation on $C$). However,
\[\sum A e^{-2\mu_A} \langle \nabla_X A X_C, X_A \rangle = \sum A e^{-2\mu_A} \langle \lambda_D^{AC} X_D + \nabla_X A X_A, X_A \rangle = \sum A X_C (\mu_A) + \lambda_A^{AC},\]
where $[X_B, X_C] = \lambda_B^{AC} X_A$. To conclude
\[-\kappa(\nabla_X A Y^A, X_C) = \ell_C X_C (\sum A \mu_A) + \ell_C \lambda_A^{AC}\]
(no summation on $C$). Next,
\[-\kappa(Y^A, \nabla_X A X_C) = -\sum A \ell_A Y^A (\nabla_X A X_C) = -\sum A \ell_A e^{-2\mu_A} \langle \nabla_X A X_C, X_A \rangle\]
\[= -\sum A \ell_A \lambda_A^{AC} - \sum A \ell_A X_C (\mu_A).\]
The momentum constraint can thus be rewritten
\[X_C (\ell_C) + \ell_C \lambda_A^{AC} - \sum A \ell_A \lambda_A^{AC} + (\ell_C - 1) X_C (\ln \theta)\]
\[+ \ell_C X_C (\sum A \mu_A) - \sum A \ell_A X_C (\mu_A) = \hat{U}(\phi) X_C (\phi)\]
(no summation on $C$). The first three terms on the left hand side have finite limits. However, the same is not in general true of the last three terms on the left hand side, nor is it in general true of the right hand side. However, it is of interest to consider
\[(117) \quad \ell_C X_C (\sum A \mu_A + \ln \theta).\]
Due to (88) and the fact that the sum of the $p_A$ equals 1, it is clear that the expression (117) converges to zero exponentially. Next, consider

$$
\sum_A \ell_A X C(\bar{\mu}_A) + \hat{U}(\phi) X C(\phi) + X C(\ln \theta).
$$

Since $\ell_A$ converges to $p_A$ exponentially, since (88) holds, since $\hat{U}(\phi)$ converges to $\Phi_\theta$ exponentially, and since $\phi + \Phi_\theta \ln \theta - \Phi_b$ converges to zero exponentially in any $C^k$-norm, this expression equals

$$
\sum_A p_A X C(-p_A \ln \theta) + \Phi_\theta X C(-\Phi_\theta + \Phi_b) + X C(\ln \theta)
$$

$$
= \left[ -\sum_A p_A X C(p_A) - \Phi_\theta X C(\Phi_\theta) \right] \ln \theta - \left[ \sum_A p_A^2 + \Phi_\theta^2 \right] X C(\ln \theta)
$$

$$
+ \Phi_\theta X C(\Phi_b) + X C(\ln \theta) = \Phi_\theta X C(\Phi_b)
$$

up to terms that decay exponentially in any $C^k$-norm. Here we used the fact that $\sum_A p_A^2 + \Phi_\theta^2 = 1$, so that

$$
\sum_A p_A X C(p_A) + \Phi_\theta X C(\Phi_\theta) = 0.
$$

Summing up, the expansion normalised limit of the momentum constraint is

$$
\mathcal{K}_C(pC) + pC\gamma^A_{AC} - \sum_A p_A \gamma^A_{AC} = \Phi_\theta \mathcal{K}_C(\Phi_b)
$$

(no summation on $C$), where $\gamma^A_{BC}$ are the structure constants associated with the frame $\{x^A\}$. On the other hand, an argument similar to the above yields

$$
(\text{div}_J_\mathcal{K})_C(x_C) = \mathcal{K}_C(pC) + pC\gamma^A_{AC} - \sum_A p_A \gamma^A_{AC}
$$

(no summation on $C$); cf. the beginning of the proof of Lemma 59. The limit of the momentum constraint can thus be written $\text{div}_J_\mathcal{K} = \Phi_\theta d\Phi_b$. \hfill \Box

5.5. Proof of Theorem 27. Finally, we are in a position to prove Theorem 27.

**Theorem 27.** The estimate (25) follows from (82) and (88). The estimate (26b) follows from (91). The estimates (26a) and (27) follow from (100). The estimate (26c) follows from (103). The existence of $\Phi_\theta$ and $\Phi_b$ as well as the conclusion that the asymptotic versions of the constraint equations hold is an immediate consequence of Lemmas 79 and 81. \hfill \Box

6. Proof of the improvements in the case of Gaussian foliations

The first step in the proof of Theorem 35 is to relate $\theta$ and $t$.

**Lemma 82.** Assume the conditions of Theorem 35 to be fulfilled. Then, using the notation of Theorem 35, there is, for each $l \in \mathbb{N}$, a constant $C_l$ such that

$$
\|\ln \theta(\cdot, t) + \ln t\|_{C^l(M_t)} + \|\theta(\cdot, t) - 1\|_{C^l(M_t)} \leq C_l t^\varepsilon
$$

for all $t \leq t_0$. In particular, (28), (14) and (16b) can be reformulated in terms of the mean curvature: there are smooth functions $b_{AB}$, $A, B = 1, \ldots, n$, and, for each $l \in \mathbb{N}$, a constant $C_l$ such that

$$
g = -dt \otimes dt + \sum_{A,B} \hat{b}_{AB} \theta^{-2p_{\min}(A,B)} \mathcal{Y} A \otimes \mathcal{Y} B,
$$

$$
\hat{h} := \sum_{A,B} \hat{b}_{AB} \mathcal{Y} A \otimes \mathcal{Y} B,
$$

$$
|D^l(\hat{h} - \hat{\bar{h}})|_{\hat{b}} \leq C_l \theta^{-\varepsilon},
$$

where the last estimate holds on $M_0$.

**Remark 83.** Given (118), the last statement of the lemma can be reversed. In other words, given (118), (119), (120) and (121), it is possible to deduce that (28), (14) and (16b) hold.

**Proof.** Due to the Hamiltonian constraint, (66) (7), p. 6], the equality (110) holds. On the other hand, the deceleration parameter $q$ is given by

$$
q = n - 1 - \frac{n^2}{n-1} \frac{\rho - \bar{\rho}}{\theta^2} - \frac{2n^2 \Lambda}{n-1} \frac{\bar{S}}{\theta^2} + n \frac{\bar{S}}{\theta^2},
$$

$$
(122)$$
where we appealed to \([66, (61), p. 27]\) and \(\bar{p} := \bar{g}^{ij} T_{ij} / n\), where \(T\) is the stress energy tensor associated with the scalar field matter. Combining \([116, 122, 60, (153), p. 48]\) and \([66, \text{Remark 84}, p. 48]\) yields
\[
q - (n - 1) = - \frac{2n}{n - 1} \frac{\Lambda}{\bar{g}} + n[(\hat{U})^2 + \text{tr}\mathcal{K}^2 - 1].
\]
Next, note that since \(\Phi_a^2 + \text{tr}\mathcal{K}^2 = 1\),
\[
(\hat{U})^2 + \text{tr}\mathcal{K}^2 - 1 = (\hat{U} \Phi - \Phi_a)^2 + 2 \Phi_a (\hat{U} \Phi - \Phi_a) + \text{tr}(\mathcal{K} - \mathcal{K})^2 + 2 \text{tr}[\mathcal{K} (\mathcal{K} - \mathcal{K})].
\]
In particular, for each \(l \in \mathbb{N}\), there is a constant \(C_l\) such that
\[
\|[(\hat{U})^2 + \text{tr}\mathcal{K}^2 - 1](\cdot, t)\|_{C^l(\bar{M})} \leq C_l t^\varepsilon
\]
for \(t \leq t_0\), due the assumptions. Next, note that
\[
\frac{d}{ds}[(\theta + \ln \theta) \circ \gamma](s) \leq C e^{\varepsilon s},
\]
where we appealed to \([78]\) and the fact that \(\theta \circ \gamma(s) = s\). Integrating this estimate from \(s_a\) to \(s_b \leq 0\), where \(s_a \leq s_b\) and \(t[\gamma(s_a)] \leq t_0\), yields
\[
s_b + \ln \theta \circ \gamma(s_b) - s_a - \ln \theta \circ \gamma(s_a) \leq C e^{\varepsilon s_b}.
\]
Assume now that \(\gamma(s_b) = (\bar{x}, t_0)\) and \(\gamma(s_a) = (\bar{x}, t)\) with \(0 < t \leq t_0\). Then this inequality yields
\[
\ln \theta(\bar{x}, t) \geq -\theta(\bar{x}, t_0) + \theta(\bar{x}, t_0) + \ln \theta(\bar{x}, t_0) - C e^{\varepsilon (\bar{x}, t_0)}.
\]
Since this is true for any \(\bar{x} \in \bar{M}\) and since \([78]\) holds,
\[
\ln \theta(\bar{x}, t) \geq -\ln t + \ln \theta_0 - C,
\]
where \(C\) is allowed to depend on \(t_0\), but not on \(\bar{x}\). Moreover, \(\theta_- := \inf_{\bar{x} \in \bar{M}} \theta(\bar{x}, t_0)\). Since \(\theta_- > 0\) by assumption, we conclude that \(\ln \theta \geq -\ln t - C\) on \(\bar{M}_0\). Combining this observation with \([125]\) yields the conclusion that \(\theta + \ln \theta\) is bounded, and thereby the conclusion that \(\ln \theta + \ln t\) is bounded. In other words, there is a constant \(C\) such that
\[
\|\ln \theta + \theta\|_{C^0(\bar{M})} + \|\ln \theta\|_{C^0(\bar{M})} \leq C
\]
for \(t \leq t_0\).
If \(\Lambda > 0\), we assume that \(\theta \geq \theta_0\) on \(\bar{M}_0\), where \(\theta_0 \in \mathbb{R}\) satisfies \(\theta_0 > [2\Lambda/(n - 1)]^{1/2}\). Then \([125]\) can be used to deduce that \(\ln \theta\) tends to infinity as \(-\sigma \ln t\) for some \(\sigma > 0\). Inserting this information into \([125]\) again leads to the conclusion that \([126]\) holds.
Next, we wish to estimate \(E_1 \ln \theta\); cf. the notation introduced at the beginning of the proof of Lemma \([61]\). Apply, to this end, \(E_1\) to \([125]\). This yields, recalling that \(\hat{U} = \theta^{-1} \partial_t\) in the present setting,
\[
-E_1(\ln \theta + \phi) \hat{U}(\theta + \ln \theta) + E_1(\phi) \hat{U}(\theta + \ln \theta) + \hat{U}[E_1(\theta + \ln \theta)]
\]
Introducing the notation \(f := E_1(\ln \theta + \phi)\), this equality can schematically be written
\[
\hat{U}(f) = g_1 f + g_2.
\]
Here \(\|g_1(\cdot, t)\|_{C^0(\bar{M})} \leq C t^\varepsilon\), where we appealed to \([78]\), \([124]\) and the fact that the right hand side of \([125]\) is \(O(t^\varepsilon)\). Evaluating this equality on an integral curve \(\gamma\) as above yields
\[
\frac{d}{ds}(f^2 \circ \gamma + 1) \geq -2(|g_1 \circ \gamma| + |g_2 \circ \gamma|)(f^2 \circ \gamma + 1).
\]
Since \(g \circ \gamma(s) = O(e^{s \gamma})\), this estimate can be integrated to yield the conclusion that \(f\) is bounded on \(M_0\). This means that \([126]\) can be improved in that the \(C^0\)-norm can be replaced by the \(C^1\)-norm. Proceeding inductively, and applying successively higher order derivatives to \([127]\), similar arguments yield the conclusion that for every \(l\), there is a constant \(C_l\) such that

\[
\|\ln \theta + \rho(\cdot, t)\|_{C^l(M)} + \|\ln \theta(\cdot, t) + \ln t\|_{C^l(M)} \leq C
\]

for all \(t \leq t_0\).

Next, consider \(\partial_t \theta^{-1} = -\theta^{-2} \partial_t \theta = (q + 1)/n\). Integrating this equality from 0 to \(t\), keeping in mind that \(\theta^{-1}\) vanishes at \(t = 0\), yields

\[
[\theta(\vec{x}, t)]^{-1} = t + \int_0^t [q(\vec{x}, s) - (n - 1)]/n \, ds
\]

Due to \([123]\), \([128]\) and the assumptions, this equality can be differentiated \(l\) times in order to deduce that

\[
\|\theta^{-1}(\cdot, t) - t\|_{C^l(M)} \leq Ct^{1 + \varepsilon}
\]

for \(t \leq t_0\). The lemma follows. \(\square\)

6.1. Proof of Theorem 35

Finally, we are in a position to prove Theorem 35.

**Theorem 35.** The idea is to verify that the conditions of Theorem 27 are satisfied and to combine the conclusions of this theorem with an estimate relating \(t\) and \(\theta\) in order to obtain the desired conclusions. Due to Lemma 82, we know that \([118]\) holds and that the metric can be represented by \([119]\). This means, in particular, that the metric can be represented as in \([22]\) with \(N = 1\). Due to \([118]\), \(\theta\) diverges to \(\infty\) uniformly. Moreover, due to \([16a]\), \([121]\) and \([118]\), it follows that \([18a]\) and \([18b]\) hold. Next, due to \([118]\) and the fact that \(N = \theta\), it follows that \([23]\) holds and that the first term on the left hand side of \([24]\) is bounded by a third times the right hand side. Since the third term on the left hand side vanishes identically (since \(N = 1\)), it remains to estimate the middle term. On the other hand,

\[
t_0 \partial_t (\phi - \Phi_a \ln t) = t_0 \partial_t \hat{U}(\phi) - \Phi_a = t_0 (\hat{U}(\phi) - \Phi_a) + (t_0 - 1) \Phi_a
\]

The first term on the right hand side is bounded by \(C_k t^\varepsilon\) in \(C^k\) and the second term on the right hand side is bounded by \(C_k t^\varepsilon\) in \(C^k\) due to \([118]\). Integrating the corresponding estimate results in the existence of a \(\Phi_b \in C^\infty(M)\) such that

\[
\|\phi(\cdot, t) - \Phi_a \ln t - \Phi_b\|_{C^l(M)} \leq C_k t^\varepsilon
\]

for all \(t \leq t_0\). This estimate implies that the second term on the right hand side of \([24]\) satisfies the desired bound. To conclude, all the conditions of Theorem 27 are satisfied. This means that all the conclusions hold, and due to \([118]\), the mean curvature \(\theta\) can be replaced by \(t^{-1}\) in the conclusions. This yields the conclusions of Theorem 35. \(\square\)

7. Obtaining data on the singularity from convergent solutions

For the remainder of the article, we assume that we have a solution which converges in the sense described at the beginning of Subsection 1.7. The goal is then to conclude that such solutions yield data on the singularity. We begin by deducing information concerning the asymptotics of the expansion normalised Weingarten map in the \(C^0\)-setting.

7.1. Asymptotics in \(C^0\). Let us begin by deriving limits for \(K\) and \(\hat{U}\).

**Lemma 84.** Let \((M, g, \phi)\) be a solution to the Einstein-scalar field equations with a cosmological constant \(\Lambda\) and a uniform volume singularity at \(t = 0\); cf. Definition 39. If \([\hat{M}]\) holds with \(k = 0\), there is a continuous \((1, 1)\)-tensor field \(\mathcal{K}\) on \(\hat{M}\) and a constant \(C\), depending only on \(C_0\), \(a_0\), \(n\) and \(\varepsilon_0\), such that

\[
|\mathcal{K} - \mathcal{K}|_{\hat{g}_{\text{ref}}} \leq C(\phi)^{a_0} e^{2c \phi}
\]

(129)
on $M$. If, in addition, \((32)\) holds, there is, analogously, a continuous function $\Phi_a$ on $\bar{M}$ such that
\[
|\bar{U}_\phi - \Phi_a| \leq C(\bar{q})^{a_0} e^{2rq}
\]
on $M$.

**Proof.** Let, to begin with, \(\{E_i\}\) be a local orthonormal frame on $T\bar{M}$ with respect to $\bar{g}_{\text{ref}}$, and let \(\{\omega^i\}\) be the dual frame. Let, moreover, $K_i^j := K(\omega^i, E_j)$. Due to \(65\) (A.4), p. 204] and \(31\),
\[
|\bar{U}_\phi| = \|L_VK(\omega^i, E_j)| \leq C_0(\bar{q})^{a_0} e^{2rq}
\]
on $M$. Let $\gamma$ be an integral curve with respect to $\bar{U}$. Then, due to \(63\) (7.9), p. 74],
\[
\frac{d}{ds}(\varrho \circ \gamma)(s) = [\bar{U}(\varrho)] \circ \gamma(s) = 1.
\]
In particular, we can parametrise $\gamma$ so that $\varrho \circ \gamma(s) = s$. Combining this observation with \(131\) yields
\[
\left| \frac{d}{ds}(K^j_i \circ \gamma)(s) \right| \leq C_0(s)^{a_0} e^{2r_2s}
\]
for all $s$ in the domain of definition of $\gamma$: note that the $\bar{M}$-component of $\gamma$ is independent of $s$, so that $\zeta$ is evaluated at a fixed point of $\bar{M}$. In particular, $K \circ \gamma$ converges to a limit and there is a $\mathcal{X}$ and a constant $C$, depending only on $C_0$, $a_0$, and $\zeta$, such that \(129\) holds. Since $\varrho$ converges to $-\infty$ uniformly, \(129\) implies that $K(\cdot, t)$ converges uniformly to $\mathcal{X}$, so that $\mathcal{X}$ has to be continuous.

The proof of \(130\) is similar. \hfill \Box

Let $p_A$ denote the eigenvalues of $\mathcal{X}$. Since the eigenvalues of $K$ are real, and the $p_A$ are limits of the eigenvalues of $K$, the $p_A$ are real. In what follows, we assume the $p_A$ to be distinct and order them so that $p_1 < \cdots < p_n$. Then
\[
|\ell_A - p_A| \leq C(\bar{q})^{a_0} e^{2rq}
\]
on $M$, where $\ell_A$ denotes the eigenvalues of $K$. Since the convergence is uniform, we can restrict the interval $I$ in such a way that the $\ell_A$ are distinct on $M$. In fact, the following holds.

**Lemma 85.** Let $(M, g, \phi)$ be a solution to the Einstein-scalar field equations with a cosmological constant $\Lambda$ and a uniform volume singularity at $t = 0$. Assume that \(31\) holds with $k = 0$. Assume, moreover, that the eigenvalues of $\mathcal{X}$ (whose existence is ensured by Lemma \(84\)) are distinct. Then, by restricting the interval $I$, if necessary, it can be assumed that the eigenvalues $\ell_A$ of $K$ are distinct on $M$ and that there is a constant $\epsilon_{\text{ref}} > 0$ such that the difference between the eigenvalues of $K$ is bounded from below by $\epsilon_{\text{ref}}$ on $M$. Moreover, by taking a finite covering space of $M$, if necessary, there is, for each $\ell_A$, a corresponding global eigenvector field $X_A$ of $K$ with $|X_A|_{\text{ref}} = 1$. Let $\{Y^A\}$ be the frame dual to $\{X_A\}$. Then there is a frame $\{X_A\}$ of continuous eigenvector fields of $\mathcal{X}$, with dual frame $\{Y^A\}$ and a constant $C$ such that
\[
|X_A - X_A|_{\text{ref}} + |Y^A - Y^A|_{\text{ref}} \leq C(\bar{q})^{a_0} e^{2r_2q}
\]
on $M$.

**Remark 86.** In what follows, we assume the conditions of this lemma to be satisfied. Moreover, we also restrict $I$ as in the lemma, and replace $\bar{M}$ with an appropriate finite cover. We also order the $\ell_A$ so that $\ell_1 < \cdots < \ell_n$. Note also that the $X_A$ are uniquely defined up to a sign (we here assume $\bar{M}$ to be connected).

**Proof.** Since $\bar{M}$ is compact and the $p_A$ are continuous functions on $\bar{M}$, there is an $\epsilon_{\text{ref}} > 0$ such that $\min_{A \neq B} |p_A - p_B| \geq 2 \epsilon_{\text{ref}}$ on $M$. Since \(132\) holds, where the convergence is uniform, we can ensure that $|\ell_A(\cdot, t) - p_A| < \epsilon_{\text{ref}}/2$ for $t \in I$ by restricting $I$, if necessary. Then $\min_{A \neq B} |\ell_A - \ell_B| \geq \epsilon_{\text{ref}}$ on $M$. Appealing to \(63\) Lemma A.1, p. 201], we can, by taking a finite covering space of $M$, if necessary, assume that there is a global frame, say $\{X_A\}$, of eigenvector fields of $K$, with
\[ \mathcal{K}X_A = \ell_A X_A \text{ (no summation), and } |X_A|_{\beta_{\text{ext}}} = 1. \] The remaining statement of the lemma is an immediate consequence of the previous observations and the fact that \((129)\) holds. \(\square\)

This lemma justifies the introduction of the standard assumptions; cf. Definition \((40)\). Next, define \(\mathcal{W}_B^A\) by the condition that
\[ (134) \hat{L}_U X_A = \mathcal{W}_B^A X_B + \mathcal{W}_A^0 U; \]
cf. [65] (6.5), p. 68. Here \(\hat{L}_U = \theta^{-1}L_U\).

**Lemma 87.** Assume the standard assumptions to hold, cf. Definition \((40)\). Define \(\bar{\mu}_A\) by the condition that \(|X_A|_{\beta} = e^{\bar{p}_A}\). Then there is a constant \(C\) and a continuous function \(r_A\) on \(M\) such that
\[ (135) |\bar{\mu}_A - p_{A\theta} - r_A| \leq C(\bar{q})^{a_0} e^{2\varrho} \]
on \(M\).

**Proof.** Combining \((31)\) for \(k = 0\) with \([65]\) (6.7)-(6.8), the fact that \(\min_{A \neq B} |\ell_A - \ell_B| \geq \ell_{\text{end}}\) on \(M\), the fact that \(|X_A|_{\beta_{\text{ext}}} = 1\) and the fact that \(|Y_A|_{\beta_{\text{ext}}}\) is bounded (this follows from \([65]\) Lemma 5.5, p. 60) implies that \(|\mathcal{W}_B^A| \leq C(\bar{q})^{a_0} e^{2\varrho} \) on \(M\). Combining this observation with \([65]\) (7.8)-(7.9), p. 74] yields
\[ (136) \hat{U}(\bar{\mu}_A - p_{A\theta}) = \ell_A + \mathcal{W}_A^A - p_A = O((\bar{q})^{a_0} e^{2\varrho}) \]
on \(M\), where we appealed to \((132)\). At this point, we can argue as in the proof of Lemma \(84\) in order to conclude that \(\bar{\mu}_A - p_{A\theta}\) converges exponentially to a continuous function \(r_A\) on \(M\). The lemma follows. \(\square\)

Next, we estimate \(\mathcal{W}_B^A\).

**Lemma 88.** Assume the standard assumptions to be satisfied and let \(\mathcal{W}_B^A\) be defined by \((134)\). Then there is a constant \(C\) such that
\[ (137) |\mathcal{W}_B^A| \leq C(\bar{q})^{a_0} e^{2\varrho} \min\{1, e^{2(p_B-p_A)\bar{e}}\} \]
on \(M\).

**Proof.** By the argument presented at the beginning of the proof of Lemma \(87\), \(|\mathcal{W}_B^A| \leq C(\bar{q})^{a_0} e^{2\varrho}\) on \(M\). Next, for \(A \neq B\), the right hand side of \([65]\) (6.19), p. 70] is antisymmetric. Thus
\[ (138) |\mathcal{W}_B^A| \leq e^{-2(\bar{\mu}_A - \bar{\mu}_B)}|\mathcal{W}_B^B| \leq C(\bar{q})^{a_0} e^{2(p_B-p_A)\bar{e}} \]
on \(M\), where we appealed to \((135)\). The lemma follows. \(\square\)

Next, let \(X_B^B\) and \(Y_B^B\) be defined as in Definition \((42)\). It is of interest to estimate the asymptotic behaviour of \(X_B^B\) and \(Y_B^A\). Due to \((133)\), we already know that there is a constant \(C\) such that
\[ (139) |X_B^B - \delta_B^B| + |Y_B^B - \delta_B^B| \leq C(\bar{q})^{a_0} e^{2\varrho} \]
on \(M\). However, these estimates can be improved.

**Lemma 89.** Assume the standard assumptions to be satisfied. Then there are constants \(a\) and \(C\) such that
\[ (140) |X_B^B - \delta_B^B| + |Y_B^B - \delta_B^B| \leq C(\bar{q})^{a_0} e^{2\varrho} \min\{1, e^{2(p_B-p_A)\bar{e}}\} \]
on \(M\).

**Proof.** Note, to begin with, that \((134)\) yields
\[ (141) \hat{U}(X_B^B) \varpi_B + X_A(\ln N)\hat{U} = \mathcal{W}_B^B X_B^B \varpi_B + \mathcal{W}_A^0 U. \]
This equality requires a comment, since the vector fields \(\varpi_B\) can, at this stage, only be assumed to be continuous. This means that \(X_B^B\) can only be assumed to be continuous. On the other hand,
the $X_B^A$ are smooth considered as functions of $t$. For this reason, the time derivative of $X_B^A$ is well defined.

One particular consequence of (141) is that $\dot{U}(X_C^A) = W_A^B X_C^B$. The goal is to use this equation in order to improve our knowledge concerning the asymptotic behaviour of the $X_A$. As a first step, note that (139) holds. Let $\gamma$ be an integral curve of $\dot{U}$ with $\varrho \circ \gamma(s) = s$, cf. the proof of Lemma 83. Then, if $A \geq 2$,

$$
\frac{d}{ds} X_A^1 \circ \gamma = W_A^1 \circ \gamma \cdot X_1^1 \circ \gamma + \sum_{B \geq 2} W_A^B \circ \gamma \cdot X_B^1 \circ \gamma.
$$

The corresponding $n - 1$ equations can be written

$$
\dot{\xi} = M \xi + \zeta,
$$

where $\xi$ is the vector consisting of the $n - 1$ elements $X_A^1 \circ \gamma$, $A = 2, \ldots, n$; $\zeta$ is the vector consisting of the $n - 1$ elements $W_A^1 \circ \gamma \cdot X_1^1 \circ \gamma$, $A = 2, \ldots, n$; $\|M(s)\| \leq C(s)^{a_0}e^{2\tau s}$; $|\xi(s)| \leq C(s)^{a_0}e^{2\tau s}$; and $|\zeta(s)| \leq C(s)^{a_0}e^{2\tau s}$; these estimates are immediate consequences of (137) and (139), and the constants (in these estimates and the ones below) are independent of $\gamma$. Then, for $-\infty < s_a \leq s \leq 0$ in the domain of definition of $\gamma$,

$$
\xi(s) = \Phi(s; s_a)\xi(s_a) + \int_{s_a}^s \Phi(s; \tau)\zeta(\tau) d\tau,
$$

where $\Phi(s; s_a)$ is defined by $\dot{\Phi}(s; s_a) = M(s)\Phi(s; s_a)$ and $\Phi(s; s_a) = \mathrm{Id}$. Due to the bounds on $\|M(s)\|$, there is a constant $C$ such that $\|\Phi(s; s_a)\| \leq C$ for all $-\infty < s_a \leq s \leq 0$ in the domain of definition of $\gamma$. This implies, in particular, that

$$
|\xi(s)| \leq C|\xi(s_a)| + C \int_{s_a}^s |\zeta(\tau)| d\tau.
$$

Letting $s_a \to -\infty$ in this estimate, keeping in mind that $|\zeta(s)| \leq C(s)^{a_0}e^{2(2^{p-1}+\epsilon)s}$, yields

$$
|\xi(s)| \leq C(s)^{a_0}e^{2(2^{p-1}+\epsilon)s}.
$$

Next, assume inductively that for some $2 \leq m \leq n - 1$, there are constants $a$ and $C$ such that

$$
|X_A^1| \leq C_2(s)^{a_0}e^{2\tau s} \max\{e^{2(p_m-p_1)\epsilon}, e^{2(p_m-p_1)\epsilon}\}
$$

for all $A \geq 2$. Note that we know this inductive assumption to hold with $m = 2$. Then, if $A > m$,

$$
\frac{d}{ds} X_A^1 \circ \gamma = \sum_{B=1}^m W_A^B \circ \gamma \cdot X_B^1 \circ \gamma + \sum_{B>m} W_A^B \circ \gamma \cdot X_B^1 \circ \gamma.
$$

Due to (137) and (144), the first sum on the right hand side is $O(s)^{a_0}e^{2(2(p_m-p_1)\epsilon)}$. Viewing (143) as an equation for $X_A^1 \circ \gamma$ for $A > m$, we can argue as above in order to conclude that the inductive hypothesis holds with $m$ replaced by $m + 1$. To conclude, there are constants $a$ and $C$ such that

$$
|X_A^1| \leq C(s)^{a}e^{2\tau s}e^{2(2(p_m-p_1)\epsilon)}
$$

for all $A \geq 2$.

As a next inductive step, assume that for some $1 \leq m \leq n - 2$, there are constants $a$ and $C$ such that

$$
|X_A^1| \leq C(s)^{a_0}e^{2\tau s}e^{2(2(2^2-m)\epsilon)}
$$

for all $B > A$ such that $A \leq m$. We know this to be true for $m = 1$. Next, for $A > m + 1$,

$$
\frac{d}{ds} X_A^{m+1} \circ \gamma = \sum_{B=m+1}^m W_A^B \circ \gamma \cdot X_B^{m+1} \circ \gamma + \sum_{B>m+1} W_A^B \circ \gamma \cdot X_B^{m+1} \circ \gamma.
$$

The first sum on the right hand side is $O((s)^{a_0}e^{2(2(p_m-p_{m+1})\epsilon)})$. We can thus argue as before in order to conclude that

$$
|X_A^{m+1}| \leq C(s)^{a_0}e^{2(2(p_{m+1}+\epsilon)\epsilon)}
$$

on $M$ for all $A > m + 1$. Assume now, inductively, that there is a $k$, satisfying $m + 1 < k \leq n - 1$, and constants $a$ and $C$ such that

$$
|X_A^k| \leq C(s)^{a}\max\{e^{2(p_k-p_{m+1})\epsilon}, e^{2(p_A-p_{m+1})\epsilon}\}
$$

for all $A > m + 1$.
on $M$ for all $A > m + 1$. We know this to be true for $k = m + 2$. If $A > k$,
\[
\frac{d}{ds} X_A^{m+1} \circ \gamma = \sum_{B \leq k} W_A^B \circ \gamma \cdot X_B^{m+1} \circ \gamma + \sum_{B > k} W_A^B \circ \gamma \cdot X_B^{m+1} \circ \gamma.
\]
Moreover, the first term on the right hand side is $O((s)^a e^{2(s-\rho_0+1)})$ and we can argue as before in order to prove that $k$ can be replaced by $k + 1$ in the inductive hypothesis (148). This means that we can replace $m$ by $m + 1$ in (147). Thus (147) holds for all $B > A$. In particular, the first term on the left hand side of (140) is bounded by the right hand side.

Finally, note that $Y_B^A$ are the components of the inverse of the matrix with components $X_B^A$. Combining this observation with the estimates for $X_B^A$ and an argument similar to the one presented at the end of the proof of Lemma 64 yields the conclusion that the second term on the left hand side of (140) satisfies the desired bound. $\square$

At this stage, we can derive asymptotics for the metric, $K$ and $\hat{\mathcal{L}}_U K$. In other words, we can prove Theorem 43.

Theorem 43. Note, to begin with, that
\[
g = -N^2 dt \otimes dt + \sum_A e^{2\tilde{\mu}_A Y^A} \otimes Y^A.
\]
In particular,
\[
g(\mathcal{Y}_B, \mathcal{X}_C) = \sum_A e^{2\tilde{\mu}_A Y^A}(\mathcal{Y}_B^A)Y^A(\mathcal{X}_C^A) = \sum_A e^{2\tilde{\mu}_A Y^A B^A C^A}.
\]
Combining this observation with (135) and (140) yields, if $A \neq B$,
\[
e^{-2p_\lambda e^{-2r_\gamma} g(\mathcal{Y}_A, \mathcal{Y}_B)} \leq C(\varrho)^a e^{2\psi_{\mu}}
\]
on $M$ (no summation). Moreover, if $r_A$ are the functions whose existence is guaranteed by Lemma 87,
\[
|e^{-2p_\lambda e^{-2r_\gamma} g(\mathcal{Y}_A, \mathcal{Y}_A)} - 1| \leq C(\varrho)^a e^{2\psi_{\mu}}
\]
on $M$ (no summation). Due to these observations, it is clear that (33) and (34) hold. That (33a) and (33b) hold follows immediately from (135) and (140) respectively.

Next, note that
\[
\mathcal{K} = \sum_A \ell_A X_A \otimes Y^A.
\]
Thus
\[
\mathcal{K}(\mathcal{Y}_B, \mathcal{X}_B) = \sum_C \ell_C \mathcal{Y}_C^A(X_C)X^A_C(\mathcal{X}_B^A) = \sum_C \ell_C X_C^A Y_C^B.
\]
Combining this equality with (132) and (140) yields the conclusion that (35c) holds. Finally, note that
\[
(\mathcal{L}_U \mathcal{K})(\mathcal{Y}_A, \mathcal{Y}_B) = (\mathcal{L}_U \mathcal{K})(X_C^A Y_C^B Y_B^D) = X_C^A (\mathcal{L}_U \mathcal{K})(Y_C^B X_B^D) Y_B^D.
\]
If $A \geq B$, then (35d) is an immediate consequence of the fact that (31) holds with $k = 0$. The interesting case is thus $A < B$. If $C = D$, then the estimate of the far right hand side of (152) is similar to the proof of (35c) for $A < B$. If $C \neq D$, the term that needs to be estimated is
\[
(\ell_D - \ell_C)X_C^A Y_C^D Y_B^D,
\]
where we appealed to (65) (6.8), p. 68]. Combining this observation with (137) and (140) yields (35d). The lemma follows. $\square$

Next, we derive conclusions concerning the mean curvature, assuming (32) to hold for $k = 0$ and the limit, say $\Phi_\alpha$, of $\hat{U} \phi$ to satisfy $\Phi_\alpha^2 + \text{tr} \mathcal{H}^2 = 1$.

Lemma 90. Assume the standard assumptions to be satisfied. Assume, moreover, that (32) holds for $k = 0$ and that the limit, say $\Phi_\alpha$, of $\hat{U} \phi$ satisfies $\Phi_\alpha^2 + \text{tr} \mathcal{H}^2 = 1$. Then there are constants $a$ and $C$ and a continuous function $r_\theta$ such that
\[
|\ln \theta + \varrho - r_\theta| \leq C(\varrho)^a e^{2\psi_{\mu}}
\]
on $M$. 

\[INITIAL DATA ON BIG BANG SINGULARITIES 45\]
Proof. Since the proof is very similar to the beginning of the proof of Lemma 82, combined with arguments similar to the proof of Lemma 84, we leave the details to the reader. \qed

7.2. Higher order derivatives. In the previous subsection, we derive conclusions in $C^0$. In the present subsection, we derive estimates in $C^k$ for some $k \geq 1$. In order to be able to do so, we need to impose additional assumptions; cf. Definition 45.

Lemma 91. Let $1 \leq k \in \mathbb{N}$ and assume the $k$-standard assumptions to hold; cf. Definition 45. Then the tensor field $\mathcal{X}$ whose existence is guaranteed by Lemma 84 is $C^k$. Moreover, there are constants $a_k$ and $C_k$ such that

\[ \sum_{l=0}^{k} | \tilde{D}^l (\mathcal{K} - \mathcal{X}) |_{\bar{g}_{\text{ref}}} \leq C_k (\notag) a_k e^{2 \varepsilon \varrho} \]

on $M$. In particular, if $\ell_A$, $X_A$, $Y^A$, $\mathcal{X}_A$ and $\mathcal{Y}^A$ are given by Definition 85 and $p_A$ are the eigenvalues of $\mathcal{X}$, ordered by size, then $\mathcal{X}_A$, $\mathcal{Y}^A$ and $p_A$ are $C^k$ and there are constants $a_k$ and $C_k$ such that

\[ | \tilde{D}^l (\ell_A - p_A) |_{\bar{g}_{\text{ref}}} + | \tilde{D}^l (X_A - \mathcal{X}_A) |_{\bar{g}_{\text{ref}}} + | \tilde{D}^l (Y^A - \mathcal{Y}^A) |_{\bar{g}_{\text{ref}}} \leq C_k (\notag) a_k e^{2 \varepsilon \varrho} \]

on $M$ for all $0 \leq l \leq k$ and all $A \in \{1, \ldots, n\}$.

Proof. Let $\{E_i\}$ be a global orthonormal frame with respect to $\bar{g}_{\text{ref}}$ and let $\{\omega^i\}$ be the dual frame. If $I = (i_1, \ldots, i_m)$, where $i_j \in \{1, \ldots, n\}$, then we use the notation $D_I := D_{E_{i_1}} \cdots D_{E_{i_m}}$, $E_I := E_{i_1} \cdots E_{i_m}$ and $|I| = m$. One can either check directly or verify by appealing to, e.g., [65, Lemma 5.7, p. 62], that if $|I| \leq k$, then

\[ | \tilde{D}_{I} \tilde{L}_U \mathcal{K} |_{\bar{g}_{\text{ref}}} \leq C_k (\notag) a_k e^{2 \varepsilon \varrho} \]

on $M$, where we appealed to (31). Next, we wish to prove that for $|I| \leq k$,

\[ |E_I [\hat{U} (\mathcal{K}^I_j)] | \leq C_k (\notag) a_k e^{2 \varepsilon \varrho} \]

on $M$, where $\mathcal{K}^I_j = \mathcal{K} (\omega^i, E_j)$. Note, to this end, that, due to [65, (A.4), p. 204],

\[ \hat{U} (\mathcal{K}^I_j) = (\hat{L}_U \mathcal{K}) (\omega^i, E_j). \]

It can therefore be demonstrated inductively that $E_I [\hat{U} (\mathcal{K}^I_j)]$ consists of a linear combination of terms of the form

\[ (\tilde{D}_{|J|} \tilde{L}_U \mathcal{K}) (\bar{D}_{|K|} \omega^I, \bar{D}_{|L|} E_j), \]

where $|J| + |K| + |L| = |I|$. Combining this observation with (155) yields the conclusion that (156) holds. Next, we wish to commute $E_I$ and $\hat{U}$. Note, to this end, that $[E_I, \hat{U}]$ can be written as a linear combination of terms of the form

\[ E_I (\ln \tilde{N}) \cdots E_{I_m} (\ln \tilde{N}) E_4 \hat{U}, \]

where $|I_1| + \cdots + |I_m| + |J| = |I|$, $|I_j| \neq 0$ and $|J| < |I|$. Combining this observation with (156) and the assumptions yields the conclusion that if $|I| \leq k$, then

\[ |\hat{U} (E_I (\mathcal{K}^I_j))| \leq C_k (\notag) a_k e^{2 \varepsilon \varrho} \]

on $M$. At this stage, we can argue as in the proof of Lemma 84 in order to conclude that $E_I (\mathcal{K}^I_j)$ converges at the rate $(\notag) a_k e^{2 \varepsilon \varrho}$. Since the convergence is uniform, we conclude that the $\mathcal{X}$ whose existence is guaranteed by Lemma 84 is in fact $C^k$. Moreover, we obtain (153). The estimate (154) is an immediate consequence of this. \qed

Next, we derive information concerning the asymptotics of the scalar field, given estimates of the form (82).
Lemma 92. Assume that the standard assumptions hold; cf. Definition 40. Assume, moreover, that (32) holds for some \( k \in \mathbb{N} \). If \( k \geq 1 \), assume, in addition, that (36) holds. Then there are \( C^k \)-functions \( \Phi_a \) and \( \Phi_b \) on \( M \) and constants \( a_k \) and \( C_k \) such that

\[
|\bar{D}^l(\hat{U}(\phi) - \Phi_a)|_{\tilde{\gamma}\ref} + |\bar{D}^l(\phi - \Phi_a \theta - \Phi_b)|_{\tilde{\gamma}\ref} \leq C_k(\theta)^{a_k}e^{2\rho}
\]
on \( M \) for all \( 0 \leq l \leq k \).

Remark 93. In order to derive conclusions concerning the scalar field, another possibility is to impose assumptions on the geometry and to deduce conclusions concerning the scalar field; cf. the proof of Lemma 91. We encourage the interested reader to provide the details.

Proof. The proof of the fact that \( \hat{U}(\phi) \) converges to a \( C^k \)-function \( \Phi_a \) and the fact that the first term on the left hand side of (158) is bounded by the right hand side is very similar to the proof of Lemma 91. We therefore omit the details. On the other hand, since (114) holds, the proof of the existence of \( \Phi_b \) and the fact that the second term on the left hand side of (158) is bounded by the right hand side is identical. The lemma follows. \( \square \)

Next, we prove that the functions \( r_A \) whose existence is guaranteed by Lemma 87 are \( C^k \).

Lemma 94. Let \( 1 \leq k \leq \mathbb{N} \) and assume the \( k \)-standard assumptions to hold; cf. Definition 15. Then the functions \( r_A \) whose existence is guaranteed by Lemma 87 are \( C^k \). Moreover, there are constants \( a_k \) and \( C_k \) such that

\[
|\bar{D}^l(\hat{\mu}_A - p_A \theta - r_A)|_{\tilde{\gamma}\ref} \leq C_k(\theta)^{a_k}e^{2\rho}
\]
on \( M \) for all \( A \in \{1, \ldots, n\} \) and \( 0 \leq l \leq k \).

Proof. Consider, to begin with, (65) (6.7)–(6.8), p. 68]. Combining these formulae with the conclusions of Lemma 91 and the assumptions yields the conclusion that

\[
|\bar{D}^l(W_A^2)|_{\tilde{\gamma}\ref} \leq C_k(\theta)^{a_k}e^{2\rho}
\]
on \( M \) for \( 0 \leq l \leq k \) and \( A, B \in \{1, \ldots, n\} \). Combining this with the first equality in (136) and (154) yields the conclusion that

\[
|E_1\hat{U}(\hat{\mu}_A - p_A \theta)| \leq C_k(\theta)^{a_k}e^{2\rho}
\]
on \( M \) for \( 0 \leq l \leq k \) and \( A \in \{1, \ldots, n\} \). Due to the observations made in the proof of Lemma 91 we can commute \( E_1 \) and \( \hat{U} \). Integrating the resulting estimate as in the proof of Lemma 87 yields the conclusion of the lemma. \( \square \)

In what follows, it is of interest to estimate the spatial derivatives of \( \hat{\mu}_A \) and \( \mu_A \).

Lemma 95. Let \( 1 \leq k \leq \mathbb{N} \) and assume the \( k \)-standard assumptions to hold; cf. Definition 15. Then there are constants \( C_k \) and \( a_k \) such that

\[
|\bar{D}^l(\mu)|_{\tilde{\gamma}\ref} + |\bar{D}^l(\hat{\mu}_A)|_{\tilde{\gamma}\ref} \leq C_k(\theta)^{a_k}
\]
on \( M \) for all \( 0 \leq l \leq k \) and \( A \in \{1, \ldots, n\} \).

Assume, in addition to the above, that (32) holds and that the limit \( \Phi_a \), whose existence is guaranteed by Lemma 92, satisfies \( \Phi_a^2 + trX^2 = 1 \). Then the function \( \theta \), whose existence is guaranteed by Lemma 90, satisfies \( \theta \in C^k(M) \). Moreover, there are constants \( L_k \) and \( b_k \) such that

\[
\sum_{l=0}^{k} |\bar{D}^l(\ln \theta + \theta - r_\theta)|_{\tilde{\gamma}\ref} \leq L_k(\theta)^{b_k}e^{2\rho}
\]
on \( M \). In particular, there are constants \( M_k \) and \( c_k \) such that

\[
\sum_{l=0}^{k} |\bar{D}^l(\mu_A - (p_A - 1)\theta - r_A - r_\theta)|_{\tilde{\gamma}\ref} \leq M_k(\theta)^{c_k}e^{2\rho}
\]
on \( M \) for all \( A \in \{1, \ldots, n\} \).
Proof. Due to (7.9), p. 74, \( \hat{U}(\varrho) = 1 \). Thus, for \( I \neq 0 \),
\[
\hat{U}(E_I \varrho) = [\hat{U}, E_I] \varrho + E_I \hat{U} \varrho = [\hat{U}, E_I] \varrho.
\]
However, the right hand side consists of a sum of terms of the form \((157)\) applied to \( \varrho \), where \(|I_1| + \cdots + |I_m| + |J| = |I| \) and \(|I| \neq 0 \); cf. the proof of Lemma 91. Since \( E_I \hat{U} \varrho \) equals zero if \( J \neq 0 \) and equals 1 if \( J = 0 \), \( \hat{U}(E_I \varrho) \) does not grow faster than polynomially in \( \varrho \). In particular, if \( \gamma \) is an integral curve of \( \hat{U} \) with \( \varrho \circ \gamma(s) = s \), then
\[
\frac{d}{ds} (E_I \varrho) \circ \gamma = f_I(s),
\]
where \(|f_I(s)| \leq C_k(s)^{ak} \) on the domain of definition of \( \gamma \), assuming \(|I| \leq k \), and \( C_k \) is a constant independent of \( \gamma \). In particular, it is clear that the first term on the left hand side of \((160)\) is bounded by the right hand side. Combining this estimate with \((159)\) yields the conclusion that \((160)\) holds.

The proof of \((161)\) is similar to the proof of Lemma 82 and we omit the details. Finally, \((162)\) is an immediate consequence of \((159)\) and \((161)\), keeping the fact that \( \mu_A = \bar{\mu}_A + \ln \theta \) in mind; cf. the comments made in connection with \[65, (3.10)–(3.11), p. 29\].

At this stage, we can derive a higher order analogue of Lemma 88.

**Lemma 96.** Let \( 1 \leq k \in \mathbb{N} \) and assume the \( k \)-standard assumptions to hold; cf. Definition 15. Then there are constants \( C_k \) and \( a_k \) such that
\[
|D^l W^A_B|_{\varrho, \text{ref}} \leq C_k \langle \varrho \rangle^{ak} e^{2(\rho_B - \rho_A) \varrho} \min \{1, e^{2(\rho_B - \rho_A) \varrho} \}
\]
on \( M \) for all \( A, B \in \{1, \ldots, n\} \) and all \( 0 \leq l \leq k \).

**Proof.** Due to (65) (6.7)–(6.8), \((154)\) and \((31)\), it is clear that
\[
|D^l W^A_B|_{\varrho, \text{ref}} \leq C_k \langle \varrho \rangle^{ak} e^{2\varrho}
\]
on \( M \) for all \( A, B \in \{1, \ldots, n\} \) and all \( 0 \leq l \leq k \). Thus \((163)\) holds for \( A \geq B \). In case \( A < B \), \((165)\) holds. Differentiating this equality, keeping \((135)\), \((160)\) and \((164)\) in mind, yields \((163)\) for \( A < B \). \( \square \)

The next goal is to derive a higher order version of Lemma 89.

**Lemma 97.** Let \( 1 \leq k \in \mathbb{N} \) and assume the \( k \)-standard assumptions to hold; cf. Definition 15. Then there are constants \( C_k \) and \( a_k \) such that
\[
|\bar{D}^l (X^A_B - \delta^A_B)|_{\varrho, \text{ref}} + |\bar{D}^l (Y^A_B - \delta^A_B)|_{\varrho, \text{ref}} \leq C_k \langle \varrho \rangle^{ak} e^{2\varrho} \min \{1, e^{2(\rho_B - \rho_A) \varrho} \}
\]
on \( M \) for all \( 0 \leq l \leq k \) and all \( A, B \in \{1, \ldots, n\} \).

**Proof.** Note, to begin with, that Lemma 91 implies that \( X^B_A \) is \( C^k \) with respect to the spatial variables and infinitely differentiable with respect to \( t \). Moreover, due to the proof of Lemma 89, \( X^B_A \) satisfies the relation \( \hat{U}(X^C_A) = W^B_A X^C_B \). Introducing \( Z^A_B = X^A_B - \delta^A_B \), this relation can be written
\[
\hat{U}(Z^C_A) = W^B_A Z^C_B + W^C_A.
\]
Let \( 0 \leq m < k \) and assume, inductively, that
\[
|E_I Z^A_B| \leq C_m \langle \varrho \rangle^{am} e^{2\varrho} \min \{1, e^{2(\rho_B - \rho_A) \varrho} \}
\]
on \( M \) for all \( A, B \in \{1, \ldots, n\} \) and \(|I| \leq m \). Note that if this estimate holds, then
\[
|E_I \hat{U} Z^A_B| \leq C_m \langle \varrho \rangle^{am} e^{2\varrho} \min \{1, e^{2(\rho_B - \rho_A) \varrho} \}
\]
on \( M \) for all \( A, B \in \{1, \ldots, n\} \) and \(|I| \leq m \). This is a consequence of combining \((167)\) with \((163)\) and \((166)\).
Due to Lemma 89, we know that for $m = 0$, the inductive hypothesis (167) holds. Next, let $I$ be such that $|I| = m + 1$ and compute the equation (166) with $E_I$. This yields

$$
\hat{U}(E_I Z^A_J) = [\hat{U}, E_I] Z^A_J + E_I (\mathcal{W}^B_A Z^C_B) + E_I (\mathcal{W}^C_A).
$$

On the other hand, $[\hat{U}, E_I]$ can be written as a linear combination of terms of the form (157), where $|I_1| + \cdots + |I_m| + |J| = |I|$. Due to this fact, the assumptions, (163), (167) and (168), the equation (169) can be rewritten

$$
\hat{U}(E_I Z^A_J) = \mathcal{W}^B_A E_I Z^C_B + F^C_A,
$$

where

$$
|F^C_A| \leq C_m(q)^a e^{2t} \min \{1, e^{2(p_B - p_A) t}\}
$$
on $M$. At this stage, an argument similar to the proof of Lemma 89 yields the conclusion that (167) holds with $m$ replaced by $m + 1$. Thus (167) holds with $m = k$. Thus the first term on the left hand side of (165) is bounded by the right hand side.

Finally, since the $Y^A_B$ are the components of the inverse of the matrix with components $X^A_B$, an argument similar to the one presented at the end of the proof of Lemma 64 yields the conclusion that the second term on the left hand side of (165) satisfies the desired bound.

Finally, we are in a position to prove Theorem 46.

Theorem 46. Since $\mathcal{X}^A$ is a frame for the cotangent space for the leaves of the foliation, it is clear that the metric can be written in the form (37). Moreover,

$$b_{AB} = e^{-2p_{\text{max}}(A,B)} g(\mathcal{X}^A, \mathcal{X}^B) = \sum \epsilon e^{2\bar{p}_\epsilon - 2p_{\text{max}}(A,B)} Y^A_A Y^C_B,$$

where we appealed to (149). Combining this equality with Lemma 94 and (165) yields the conclusion that (38) holds.

The estimates (39a) and (39b) follow immediately from (159) and (165) respectively. Next, note that (150) and (151) hold. Combining these equalities with (154) and (165) yields (39c). The proof of (39d) is similar to the end of the proof of Theorem 43, given the assumptions, (154), (163) and (165).

The estimate (40) is an immediate consequence of (158) and (41) is an immediate consequence of (161). The equality (42a) and the estimate (42b) follow by combining (41) with (37) and (38). Next, the estimates (44a) and (44b) follow from (39a) and (39b) respectively. The estimate (44b) follows from (162). Next, (44c) and (44d) follow from (161), (39c) and (39d). The estimate (45) follows from (40), (123), (124), (153), (160) and (161).

Next, assume $N = 1$ and compute

$$\partial_t \theta^{-1} = -\theta^{-1} \partial_t \ln \theta = -n^{-1} \hat{U}(n \ln \theta) = n^{-1}(q + 1) = 1 + n^{-1}(q - (n - 1)).$$

Since the right hand side converges uniformly to 1 as $t \to 0+$, there is a $T > 0$ such that on $(0, T]$, $1/2 \leq \partial_t \theta \leq 2$ on $M \times (0, T]$. Since $\theta \to \infty$ uniformly as $t \to 0+$, we conclude that $t/2 \leq \theta^{-1} \leq 2t$ on $(0, T]$. Combining this observation with (170) and (45) yields the conclusion that there are constants $C_k$ and $a_k$ such that

$$|D^l(\partial_t \theta^{-1} - 1)|_{\bar{g}_{ref}} \leq C_k (\ln t)^{a_k} t^{2e}$$
on $M$ for $0 \leq l \leq k$. Integrating this estimate yields (46). Combining (46) with (42a) and (42b) yields (47a) and (47b).

7.3. Reproducing the conditions on data on the singularity. Next, we prove Theorem 49.

Theorem 49. Recall, to begin with, the Hamiltonian constraint (116). Due to Lemma 90, it is clear that $\Omega_A$ decays as $e^{2t}$. Next, due to (116), (153), p. 48,

$$2\Omega + \text{tr} \kappa^2 = (\hat{U}\phi)^2 + \text{tr} \kappa^2 + \sum_A e^{-2p_A} (X_A \phi)^2.$$
Due to (133), (158) and (160), we know that $X_A \phi$ does not grow faster than polynomially in $\theta$. Next, note that since the sum of the $p_A$ is bounded from above by $1$, since there is an $\epsilon_{nd}>0$ such that $|p_A - p_B| \geq \epsilon_{nd}$, and since $n \geq 3$, there is an $\epsilon_{sp}>0$ such that $p_A \leq 1 - \epsilon_{sp}$. Combining this observation with (162), yields the conclusion that there is a constant $C$ such that

$$\mu_A \geq -\epsilon_{sp} \theta - C. \quad (171)$$

Combining the above observations,

$$\theta^{-2} \bar{S} = (\bar{U} \phi)^2 + \text{tr} K^2 - 1 + \sum_A e^{-2\mu_A}(X_A \phi)^2 + 2\Omega_A. \quad (172)$$

Due to (129), (158), the fact that $\Phi_{\theta}(172)$ this observation with (162) yields the conclusion that there is a constant $C$ such that

$$\theta^{-2} \bar{S} < C e^{2\eta \theta}. \quad (173)$$

Next, define $\lambda_{BC}^A$, by $[X_B, X_C] = \lambda_{BC}^A X_A$. Due to (154), it is clear that $\lambda_{BC}^A$ and $X_B(\lambda_{BC}^A)$ converge in $C^\infty$. Due to (154) and (160), it is clear that $X_A(\bar{\mu}_B)$ and $X_B X_B(\bar{\mu}_C)$ do not grow faster than polynomially. On the other hand, $e^{-\mu_A}$ decays to zero exponentially due to (171).

Combining these observations with the formula [66] (199), p. 58 (with $\gamma_{BC}^A$ replaced by $\lambda_{BC}^A$) yields

$$\theta^{-2} \bar{S} + \frac{1}{4} \sum A, B, C e^{2\mu_B - 2\mu_A - 2\mu_C} (\lambda_{AC}^B)^2 \leq C(\theta)^2 e^{2\epsilon_{sp} \theta} \quad (174)$$

on $M$. Combining this estimate with (173) yields the conclusion that there is an $\eta > 0$ and a $C$ such that

$$\sum A, B, C e^{2\mu_B - 2\mu_A - 2\mu_C} (\lambda_{AC}^B)^2 \leq C e^{2\eta \theta} \quad (175)$$

on $M$. Assume now that $\bar{x} \in M$ is such that $1 + p_A(\bar{x}) - p_B(\bar{x}) - p_C(\bar{x}) \leq 0$. Then, due to the continuity of the $p_A$, there is an open neighbourhood $V$ of $\bar{x}$ such that $1 + p_A - p_B - p_C \leq \eta/2$ in $V$. On the other hand, due to (162),

$$2\mu_A - 2\mu_B - 2\mu_C = 2(1 + p_A - p_B - p_C) + O(1) \geq \eta \theta - C$$

on $V$, for some constant $C$. Combining this estimate with (175) yields

$$(\lambda_{BC}^A)^2 \leq C e^{\eta \theta} \quad (176)$$

on $V$. This means that $\lambda_{BC}^A$ converges to zero on $V$. On the other hand, $\lambda_{BC}^A$ converges to $\gamma_{BC}^A$ due to (154). This proves the first statement of the theorem; note that $A$, $B$ and $C$ have to be distinct in order for the statement not to be void, in which case $A\mathcal{Y}^A([\mathcal{X}_B, \mathcal{X}_C]) = 0$ is equivalent to $A\mathcal{Y}^A([\mathcal{X}_B, \mathcal{X}_C]) = 0$.

Next, note that $\text{tr} \mathcal{X} = 1$ is an immediate consequence of the fact that $\mathcal{K} \to \mathcal{X}$ and the fact that $\text{tr} \mathcal{K} = 1$. The fact that $\mathcal{X}$ is symmetric with respect to $\tilde{h}$ is an immediate consequence of the fact that the fields $\mathcal{X}$ is orthonormal with respect to $\tilde{h}$. The relation $\text{tr} \mathcal{X}^2 + \Phi_{\theta}^2 = 1$ holds by assumption. That (53a) holds is an immediate consequence of (153), (161), (121) and the definition of $\tilde{h}$; cf. (52). That (53b) holds is an immediate consequence of (51). What remains to be demonstrated is that $\text{div} \mathcal{X} = \Phi_{\theta} \delta \Phi_{\theta}$. The template for doing so is the proof of Lemma [51]. However, the definition of $X_A$ etc. used in this proof differs from the definition of the present section. Let $X_{\Lambda} = e^{-\theta A} X_A$ (no summation) and define $e^{\theta A}$ by $[X_{\Lambda} A] = e^{\theta A}$. Then $\{X_{\Lambda}\}$ is a $C^k$-frame with dual $C^k$-frame $\{Y^A\}$. Note also that $\{X_{\Lambda}\}$ satisfies the conditions stated in Remark [14] except for the regularity; the vector fields $\tilde{X}_{\Lambda}$ are $C^k$. Next, since $\bar{r}_A = r_A + p_A r_\theta$,

$$\bar{\mu}_A + p_A \ln \theta = \bar{\mu}_A - \bar{r}_A + p_A \ln \theta = (\mu_A - p_A \bar{\rho} - r_A) + p_A \bar{\rho} + r_A - p_A r_\theta + p_A \ln \theta = (\mu_A - p_A \bar{\rho} - r_A) + p_A (\bar{\rho} + \ln \theta - r_\theta).$$

Combining this equality with (159) and (161) yields

$$\| \bar{\theta}^\prime (\bar{\mu}_A + p_A \ln \theta) \|_{\theta r} \leq C_k (\theta) \delta c e^{2\rho \theta}$$
on $M$ for all $A$ and all $0 \leq l \leq k$. Given the above estimates, we can revisit the proof of Lemma 51 with suitable substitutions: $X_A$ and $Y^A$ are replaced by $\bar{X}_A$ and $\bar{Y}^A$ respectively; $\bar{\mu}_A$ is replaced by $\bar{\mu}_A$; $\Phi_a$ and $\Phi_b$ are replaced by $\bar{\Phi}_a$ and $\bar{\Phi}_b$ respectively; the structure constants are calculated with respect to $\bar{X}_A$ etc. The conclusion is that $\text{div}_\bar{h}\mathcal{K} = \bar{\Phi}_a d\bar{\Phi}_b$. The theorem follows. $\square$

Finally, we are in a position to prove Theorem 51.

**Theorem 51.** Note that the conditions of Theorems 46 and 49 are fulfilled. In particular, it thus follows that the metric can be represented as in (47a). The corresponding $\tilde{h}$ is given by

$$\tilde{h} := \sum_{A,B} c_{AB} Y^A \otimes Y^B.$$ 

Due to (48), it follows that the limit of $\tilde{h}$ is the $\bar{h}$ appearing in the statement of Theorem 49 in fact,

$$\sum_{i=0}^k |D^i(\tilde{h} - \bar{h})|_{\gamma_{ref}} \leq C_k (\ln t)^{a_k t^{2\epsilon}}.$$ 

Moreover, the limit of $\mathcal{K}$, i.e. $\mathcal{K}$, is of course the same. This means that $\text{div}_\bar{h}\mathcal{K} = 0$, that $\text{tr}\mathcal{K}^2 = 1$, $\text{tr}\mathcal{K} = 1$ and that $\mathcal{K}$ is symmetric with respect to $\bar{h}$. Moreover, $Y^1[\mathcal{F}_2, \mathcal{F}_3] = 0$. Assuming $\mathcal{K}$ to be $C^\infty$, Lemma 55 yields coordinates $(V, \bar{x})$ such that the conclusions of Lemma 55 hold; denote the renormalised bases by $\{\mathcal{F}_A\}$ and $\{\mathcal{Y}^A\}$. Then $\mathcal{Y}^A = \beta_A \mathcal{Y}^A$ (no summation) for some positive functions $\beta_A$. Moreover,

$$\sum_{A,B} c_{AB} t^{2p_{\max}(A,B)} Y^A \otimes Y^B = \sum_{A,B} \bar{c}_{AB} t^{2p_{\max}(A,B)} \bar{Y}^A \otimes \bar{Y}^B = \sum_{i,j} \bar{d}_{ij} t^{2p_{\max}(i,j)} d\bar{x}^i \otimes d\bar{x}^j,$$

where $\bar{c}_{AB} = \beta_A \bar{c}_{AB} \beta_B$ (no summation),

$$\sum_{i=0}^k |D^i(\bar{c}_{AB} - \bar{c}_{AB}^2 \delta_{AB})| \leq C (\ln t)^{a_k t^{2\epsilon}}$$

and $\bar{c}_A = \beta_A \bar{c}_A$ (no summation); cf. (47b). Moreover, it can be deduced that (51) holds on $W$, where $\bar{d}_{ii} = \bar{c}_{ii}^2 > 0$. Combining this observation with Lemma 55 yields (55). Next, note that (44d) can be written

$$(176) \quad \sum_{i=0}^k |D^i[\mathcal{K}(\mathcal{Y}^A, \mathcal{Y}^B) - \mathcal{K}(\mathcal{Y}^A, \mathcal{Y}^B)]| \leq C_k (\ln t)^{a_k t^{2\epsilon}} \min\{1, t^{2(p_B - p_A)}\}.$$ 

One immediate consequence of this estimate is that

$$\sum_{i=0}^k |D^i[\mathcal{K}_i - \mathcal{K}_i]|_{\gamma_{ref}} \leq C_k (\ln t)^{a_k t^{2\epsilon}}$$

on $W$ for $i \geq j$, where $\mathcal{K}_i = \mathcal{K}(d\bar{x}^i, \partial_{\bar{x}^i})$. On the other hand,

$$\mathcal{K}(d\bar{x}^3, \partial_{\bar{x}}) = \mathcal{K}(\bar{Y}^1, \bar{Y}^2 - \bar{Y}^3 \bar{\beta}_2) = O_k((\ln t)^{a_k t^{2(p_B - p_A)}})$$

on $W$, where we appealed to (176) and Lemma 55 and the notation $O_k$ signifies that the relevant estimates hold for up to $k$ derivatives. The arguments concerning the remaining components are similar and the estimate

$$(177) \quad \sum_{i=0}^k |D^i[\mathcal{K}_i - \mathcal{K}_i]|_{\gamma_{ref}} \leq C_k (\ln t)^{a_k t^{2\epsilon}} \min\{1, t^{2(p_1 - p)}\}$$

follows (on $W$). By a similar argument, appealing to (44f),

$$(178) \quad \sum_{i=0}^k |D^i[\mathcal{\hat{K}}_U \mathcal{K}]_i| \leq C_k (\ln t)^{a_k t^{2\epsilon}} \min\{1, t^{2(p_1 - p)}\}$$

on $W$. Combining these estimates with (44f), it can be deduced that (56) holds. Finally, combining the fact that $\text{div}_\bar{h}\mathcal{K} = 0$ with the arguments presented in the proof of Lemma 55 yields the final conclusion of the lemma. $\square$
References

[1] Ames, E.; Beyer, F.; Isenberg, J.; LeFloch, P. G.: Quasilinear hyperbolic Fuchsian systems and AVTD behavior in $T^d$-symmetric vacuum spacetimes. Ann. Henri Poincaré 14, no. 6, 1445–1523 (2013)
[2] Ames, E.; Beyer, F.; Isenberg, J.; LeFloch, P. G.: A class of solutions to the Einstein equations with AVTD behavior in generalized wave gauges. J. Geom. Phys. 121, 42–71 (2017)
[3] Andersson, L.; Fajman, D.: Nonlinear stability of the Milne model with matter. Comm. Math. Phys. 378, no. 1, 261–298 (2020)
[4] Andersson, L.; Moncrief, V.: Future complete vacuum spacetimes. The Einstein equations and the large scale behavior of gravitational fields, 299–330, Birkhäuser, Basel (2004)
[5] Andersson, L.; Moncrief, V.: Einstein spaces as attractors for the Einstein flow. J. Differential Geom. 89, no. 1, 1–47 (2011)
[6] Andersson, L.; Rendall, A. D.: Quiescent cosmological singularities. Comm. Math. Phys. 218, 479–511 (2001)
[7] Andréasson, H.; Ringström, H.: Proof of the cosmic no-hair conjecture in the $T^3$-Gowdy symmetric Einstein-Vlasov setting. J. Eur. Math. Soc. (JEMS) 18, no. 7, 1565–1650 (2016)
[8] Béguin, F.: Aperiodic oscillatory asymptotic behavior for some Bianchi spacetimes. Class. Quantum Grav. 27, 185005 (2010)
[9] Belinski, V. A.; Khalatnikov, I. M.; Lifshitz, E. M.: Oscillatory approach to a singular point in the relativistic cosmology. Adv. Phys. 19, 525–573 (1970)
[10] Belinski, V. A.; Khalatnikov, I. M.: Effect of scalar and vector fields on the nature of the cosmological singularity. Sov. Phys. JETP 36, 591–597 (1973)
[11] Belinski, V. A.; Khalatnikov, I. M.; Lifshitz, E. M.: A general solution of the Einstein equations with a time singularity. Adv. Phys. 31, 639–667 (1982)
[12] Bieri, L.: Part I: Solutions of the Einstein vacuum equations. Extensions of the stability theorem of the Minkowski space in general relativity, 1–295, AMS/IP Stud. Adv. Math., 45. Amer. Math. Soc., Providence, RI, (2009)
[13] Bieri, L.; Zipser, N.: Extensions of the stability theorem of the Minkowski space in general relativity. AMS/IP Studies in Advanced Mathematics, 45. Amer. Math. Soc., Providence, RI; International Press, Cambridge, MA, xxiv+491 pp. (2009)
[14] Bigorgne, L.; Fajman, D.; Joudioux, J.; Smulevici, J.; Thaller, M.: Asymptotic stability of Minkowski spacetime with non-compactly supported massless Vlasov matter. Arch. Ration. Mech. Anal. 242, no. 1, 1–147 (2021)
[15] Brehm, B.: Bianchi VIII and IX vacuum cosmologies: Almost every solution forms particle horizons and converges to the Mixmaster attractor. Preprint, arXiv:1606.08058
[16] Calogero, S.; Heinzle, J. M.: Oscillations toward the singularity of locally rotationally symmetric Bianchi type IX cosmological models with Vlasov matter. SIAM J. Appl. Dyn. Syst. 9, no. 4, 1244–1262 (2010)
[17] Calogero, S.; Heinzle, J. M.: Bianchi cosmologies with anisotropic matter: locally rotationally symmetric models. Phys. D 240, no. 7, 636–669 (2011)
[18] Christodoulou, D.; Klainerman, S.: The global nonlinear stability of the Minkowski space. Princeton Mathematical Series, 41. Princeton University Press, Princeton, NJ. x+514 pp. (1993)
[19] Dafermos, M.; Holzegel, G.; Rodnianski, I.; Taylor, M.: The non-linear stability of the Schwarzschild family of black holes. Preprint, arXiv:2104.08222
[20] Dafermos, M.; Luk, J.: The interior of dynamical vacuum black holes I: The C0-stability of the Kerr Cauchy horizon. Preprint, arXiv:1710.01722
[21] Damour, T.; de Buyl, S.: Describing general cosmological singularities in Iwasawa variables. Phys. Rev. D 77, no. 4, 043520, 26 pp. (2008)
[22] Damour, T.; Henneaux, M.; Rendall, A. D.; Weaver, M.: Kasner-like behaviour for subcritical Einstein-matter systems. Ann. Henri Poincaré 3, 1049–1111 (2002)
[23] Damour, T.; Henneaux, M.; Nicolai, H.: Cosmological billiards. Class. Quantum Grav. 20, R145–R200 (2003)
[24] Damour, T.; Nicolai, H.: Higher order M theory corrections and the Kac-Moody algebra E10. Class. Quantum Grav. 22, 2849–2879 (2005)
[25] Dullemond, T.: Chaotic dynamics of spatially homogeneous spacetimes. Université Paris 13 - Sorbonne Paris Cité, 2019. https://tel.archives-ouvertes.fr/tel-02488655
[26] Fajman, D.; Joudioux, J.; Smulevici, J.: The stability of the Minkowski space for the Einstein-Vlasov system. Anal. PDE 14, no. 2, 425–531 (2021)
[27] Fajman, D.; Oliynyk, T. A.; Wyatt, Z.: Stabilizing relativistic fluids on spacetimes with non-accelerated expansion. Comm. Math. Phys. 383, no. 1, 401–426 (2021)
[28] Fajman, D.; Wyatt, Z.: Attractors of the Einstein-Klein-Gordon system. Comm. Partial Differential Equations 46, no. 1, 1–30 (2021)
[29] Fournodavlos, G. and Luk, J.: Asymptotically Kasner-like singularities. Preprint, arXiv:2003.13591
[30] Fournodavlos, G.; Rodnianski, I.; Speck, J.: Stable Big Bang Formation for Einstein’s Equations: The Complete Sub-critical Regime. Preprint, arXiv:2012.05888
[31] Friedmann, H.: Existence and structure of past asymptotically simple solutions of Einstein’s field equations with positive cosmological constant. J. Geom. Phys. 3, no. 1, 101–117 (1986)
[32] Friedrich, H.: On the existence of n-geodesically complete or future complete solutions of Einstein’s field equations with smooth asymptotic structure. Comm. Math. Phys. 107, no. 4, 587–609 (1986)

[33] Friedrich, H.: On the global existence and the asymptotic behavior of solutions to the Einstein-Maxwell-Yang-Mills equations. J. Differential Geom. 34, no. 2, 275–345 (1991)

[34] Hadžić, M.; Speck, J.: The global future stability of the FLRW solutions to the dust-Einstein system with a positive cosmological constant. J. Hyperbolic Differ. Equ. 12, no. 1, 87–188 (2015)

[35] Heinzle, J. M.; Uggla, C.: A new proof of the Bianchi type IX attractor theorem. Class. Quantum Grav. 26, no. 7, 075015 (2009)

[36] Heinzle, J. M.; Uggla, C.; Niklas Rohr, N.: The cosmological billiard attractor. Adv. Theor. Math. Phys. 13, 293–307 (2009)

[37] Heinzle, J. M.; Uggla, C.; Lim, W. C.: Spike oscillations. Phys. Rev. D 86, 104049 (2012)

[38] Demaret, J.; Hennaux, M.; Spindel, P.: Nonoscillatory Behavior In Vacuum Kaluza-Klein Cosmologies. Phys. Lett. 164B, 27–30 (1985)

[39] Hintz, Peter; Vasy, Andras The global non-linear stability of the Kerr-de Sitter family of black holes. Acta Math. 220, no. 1, 1–206 (2018)

[40] Hintz, P.; Vasy, A.: Stability of Minkowski space and polyhomogeneity of the metric. Ann. PDE 6, no. 1, Paper No. 2, 146 pp. (2020)

[41] Isenberg, J.; Kichenassamy, S.: Asymptotic behavior in polarized $T^2$-symmetric vacuum space-times. J. Math. Phys. 40, no. 1, 340–352 (1999)

[42] Isenberg, J.; Luo, X.: Power law inflation with electromagnetism. Ann. Phys. 334, 429–454 (2013)

[43] Isenberg, J.; Moncrief, V.: Asymptotic behaviour in polarized and half-polarized $U(1)$ symmetric vacuum spacetimes. Class. Quantum Grav. 19, no. 21, 5361–5386 (2002)

[44] Klingen, S.; Rendall, A.: Analytic description of singularities in Gowdy spacetimes. Class. Quantum Grav. 15, no. 5, 1339–1355 (1998)

[45] Klainerman, S.; Szeftel, J.: Global nonlinear stability of Schwarzschild spacetime under polarized perturbations. Annals of Mathematics Studies, 210. Princeton University Press, Princeton, NJ, 2020. xi+840 pp.

[46] Klainerman, S.; Szeftel, J.: Global nonlinear stability of Schwarzschild spacetime under polarized perturbations. Comm. Pure Appl. Math. 71, no. 6, 918–976 (2018)

[47] Klainerman, S.; Szeftel, J.: Geometric analysis in general relativity. Princeton University Press, Princeton, NJ, 2012. xiv+393 pp.

[48] Klainerman, S.; Szeftel, J.: Global nonlinear stability of Schwarzschild spacetime under polarized perturbations. Ann. of Math. (2) 187, no. 2, 453–509 (2018)

[49] Klainerman, S.; Szeftel, J.: Geometric analysis in general relativity. Princeton University Press, Princeton, NJ, 2012. xiv+393 pp.

[50] Klainerman, S.; Szeftel, J.: Global nonlinear stability of Schwarzschild spacetime under polarized perturbations. Comm. Pure Appl. Math. 71, no. 6, 918–976 (2018)

[51] Klainerman, S.; Szeftel, J.: Geometric analysis in general relativity. Princeton University Press, Princeton, NJ, 2012. xiv+393 pp.

[52] Klainerman, S.; Szeftel, J.: Global nonlinear stability of Schwarzschild spacetime under polarized perturbations. Ann. of Math. (2) 187, no. 2, 453–509 (2018)

[53] Klainerman, S.; Szeftel, J.: Geometric analysis in general relativity. Princeton University Press, Princeton, NJ, 2012. xiv+393 pp.

[54] Klainerman, S.; Szeftel, J.: Global nonlinear stability of Schwarzschild spacetime under polarized perturbations. Ann. of Math. (2) 187, no. 2, 453–509 (2018)

[55] Klainerman, S.; Szeftel, J.: Geometric analysis in general relativity. Princeton University Press, Princeton, NJ, 2012. xiv+393 pp.
[70] Rodnianski, I.; Speck, J.: Stable big bang formation in near-FLRW solutions to the Einstein-scalar field and Einstein-stiff fluid systems. Selecta Math. (N.S.) 24, no. 5, 4293–4459 (2018)
[71] Rodnianski, I.; Speck, J.: On the nature of Hawking’s incompleteness for the Einstein-vacuum equations: the regime of moderately spatially anisotropic initial data. J. Eur. Math. Soc. (JEMS) 24, no. 1, 167–263 (2022)
[72] Speck, J.: The nonlinear future stability of the FLRW family of solutions to the Euler–Einstein system with a positive cosmological constant. Selecta Math. (N.S.) 18, 633–715 (2012)
[73] Speck, J.: The stabilizing effect of spacetime expansion on relativistic fluids with sharp results for the radiation equation of state. Arch. Ration. Mech. Anal. 210, 535–579 (2013)
[74] Speck, J.: The maximal development of near–FLRW data for the Einstein–scalar field system with spatial topology $S^3$. Comm. Math. Phys. 364, no. 3, 879–979 (2018)
[75] Ståhl, F.: Fuchsian analysis of $S^2 \times S^1$ and $S^3$ Gowdy spacetimes. Class. Quantum Grav. 19, no. 17, 4483–4504 (2002)
[76] Svedberg, C.: Future stability of the Einstein–Maxwell-scalar field system. Ann. Henri Poincaré 12, 849–917 (2011)
[77] Weaver, M.: Dynamics of magnetic Bianchi VIo-cosmologies. Class. Quantum Grav. 17, no. 2, 421–434 (2000)
[78] Zipser, N.: The global nonlinear stability of the trivial solution of the Einstein-Maxwell equations. Thesis (Ph.D.)–Harvard University. 198 pp. (2000)