## MEAN CURVATURE FLOW

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# 1. Mean curvature flow

Let  $M_0 \subset \mathbb{R}^{n+1}$  be a smooth n-dimensional hypersurface without boundary, given by an immersion  $X_0: M^n \to \mathbb{R}^{n+1}$ , where  $M^n$  is an abstract smooth manifold. We consider the family of embeddings  $X: M^n \times [0,T) \to \mathbb{R}^{n+1}$  with

$$X(p,0) = X_0(p)$$

for all  $p \in M^n$  and

$$\partial_t X(p,t) = \mathbf{H}(p,t) = -H(p,t)\boldsymbol{\nu}(p,t) = \Delta_{M_t} X(p,t)$$
 (MCF)

for all  $(p,t) \in M^n \times [0,T)$ . We abbreviate  $M_t := X(M^n,t)$ . In the following, we will write  $\Delta := \Delta_{M_t}$  and  $\nabla := \nabla^{M_t}$ .

**Example 1.1** (Shrinking spheres and cylinders). (i) Let  $M_t = \mathbb{S}_{r(t)}^n$ , then (MCF) reduces to an ODE for the radius, namely

$$r' = -\frac{n}{r} \,.$$

The solution with r(0) = R is

$$r(t) = \sqrt{R^2 - 2nt} \,,$$

for 
$$t \in (-\infty, R^2/2n)$$
.

- (ii) The shrinking cylinders  $M_t = \mathbb{S}^m_{r(t)} \times \mathbb{R}^{n-m}$  with  $r(t) = \sqrt{R^2 2mt}$  exist for  $t \in (-\infty, R^2/2m)$ .
- (iii) For n = 1 the so-called grim reaper is given by  $M_t = \text{graph}(u_t)$ , where  $u(x, t) = t \log \cos x$  with  $x \in (-\pi, \pi)$ .

**Remark 1.2** (Normal motion and tangential diffeomorphisms). See [Eck04, Remark 2.2(3)]. We will often consider smoothly embedded hypersurfaces  $M_t$  satisfying

$$(\partial_t x)^{\perp} = \langle \partial_t x, \boldsymbol{\nu}(x) \rangle \boldsymbol{\nu}(x) = \mathbf{H}(x)$$

for  $x \in M_t$ , where  $\perp$  denotes the projection onto the normal space of  $M_t$ . This equation is equivalent to (MCF) up to diffeomorphisms tangent to  $M_t$ . Indeed, let  $\tilde{X}(\cdot,t):M^n\to\mathbb{R}^{n+1}$  with  $M_t=\tilde{X}(M^n,t)$  be a family of embeddings satisfying the equation

$$\left(\partial_t \tilde{X}(q,t)\right)^{\perp} = \tilde{\mathbf{H}}(q,t) := \mathbf{H}\left(\tilde{X}(q,t)\right)$$

for  $q \in M^n$ , where  $\perp$  denotes the projection onto the normal space of  $\tilde{X}(M^n, t)$ . Let  $\phi_t = (\cdot, t)$  be a family of diffeomorphisms of  $M^n$  satisfying

$$\nabla \tilde{X}(\phi(p,t),t)\partial_t \phi(p,t) = -\left(\partial_t \tilde{X}(\phi(p,t),t)\right)^{\top}$$

where  $\top$  denotes projection onto the tangent space of  $\tilde{X}(M^n, t)$ . The local existence of such a family is guaranteed by the assumptions on  $\tilde{X}$ . If we set

$$X(p,t) = \tilde{X}(\phi(p,t),t)$$

then  $M_t = X(M^n, t) = \tilde{X}(M^n, t)$ , and

$$\partial_t X(p,t) = \partial_t \tilde{X}(p,t) + \nabla \tilde{X}(\phi(p,t),t) \partial_t \phi(p,t) = \left(\partial_t \tilde{X}(q,t)\right)^{\perp} = \mathbf{H}(X(p,t)).$$

The previous remark results in the following theorem (see [Sch17a, Theorem 10.6]).

**Theorem 1.3.** Let  $X: M^n \times [0,T) \to \mathbb{R}^{n+1}$  be a solution to (MCF), that is  $\langle \partial_t X, \boldsymbol{\nu} \rangle = -H$ . Let  $R \in O(n+1)$  be an orthonormal map and  $\phi: M^n \times [0,T) \to M^n$  smooth. so that  $\phi(\cdot,t)$  is a diffeomorphism. Then  $\tilde{X}(p,t) := RX(\phi(p,t),t)$  evolves by

$$\left\langle \partial_t \tilde{X}(p,t), \tilde{\boldsymbol{\nu}}(p,t) \right\rangle = -\tilde{H}(p,t),$$

where  $\tilde{H}(p,t) = H(\phi(p,t),t)$  for all  $p \in M^n$  and  $t \in [0,T)$ .

**Lemma 1.4** (Evolution equations). Let  $(M_t)_{t\in[0,T)}$  evolve by (MCF). Then,

$$\begin{split} \partial_t g_{ij} &= -2Hh_{ij}\,,\\ \partial_t g^{ij} &= 2Hh^{ij}\,,\\ \partial_t d\mu_t^n &= -H^2\,d\mu_t^n\,,\\ \partial_t \boldsymbol{\nu} &= \nabla H\,,\\ \partial_t h_{ij} &= \nabla_i \nabla_j H - Hh_i^k h_{jk}\\ &= \Delta h_{ij} - 2Hh_i^k h_{jk} + |A|^2 h_{ij}\,,\\ \partial_t h_j^i &= \Delta h_j^i + |A|^2 h_j^i\,,\\ \partial_t H &= \Delta H + H|A|^2\,,\\ \partial_t |A|^2 &= \Delta |A|^2 - |\nabla A|^2 + 2|A|^4\,,\\ \partial_t |\nabla^m A|^2 &\leq \Delta |\nabla^m A|^2 - 2|\nabla^{m+1} A|^2\\ &+ C(m,n) \sum_{i+j+k=m} |\nabla^m A| \cdot |\nabla^i A| \cdot |\nabla^j A| \cdot |\nabla^k A| \end{split}$$

for all  $t \in [0,T)$ .

Proof. See e.g. [Sch18, Section 3].

Corollary 1.5. We have that

$$\partial_t \mu_t^n(M_t) = -\int_{M_t} H^2 d\mu_t^n \,.$$

Moreover, (MCF) is the negative  $L^2$  gradient flow for the surface area functional.

*Proof.* For arbitrary normal speeds  $\partial_t X = -F \nu$ , we have that  $\partial_t g_{ij} = -2F h_{ij}$  and

$$\frac{d}{dt} \int_{M_t} d\mu_t^n = -\int_{M_t} F H d\mu_t^n \ge -\left(\int_{M_t} F^2 d\mu_t^n\right)^{1/2} \left(\int_{M_t} H^2 d\mu_t^n\right)^{1/2}$$

with equality if and only if F = H.

**Theorem 1.6** (Short time existence). Let  $M_0 \subset \mathbb{R}^{n+1}$  be a smooth, compact hypersurface given by an immersion  $X_0: M^n \to \mathbb{R}^{n+1}$ , there exists a unique, smooth solution of (MCF) in some positive time interval.

*Proof.* See e.g. [Man11, Section 1.5].  $\Box$ 

Remark 1.7. See [Man11, Remark 1.5.4]. To proof existence and uniqueness for noncompact initial surfaces one needs estimates on the initial hypersurface (like similarly, on the initial datum in order to deal with the heat equation in all  $\mathbb{R}^n$ ) to have existence in some positive interval of time. One possibility is to assume a uniform control on the norm of the second fundamental form of the initial hypersurface. Ecker and Huisken [EH89] showed that a uniform local Lipschitz condition on a hypersurface is sufficient to guarantee short time existence.

**Theorem 1.8** (Comparison principle). Let  $X:M^n\times [0,T)\to \mathbb{R}^{n+1}$  and  $Y:N^n\times [0,T)\to \mathbb{R}^{n+1}$  be two hypersurfaces moving by MCF, where  $M^n$  is compact. Then the distance between them is nondecreasing in time.

*Proof.* We follow the lines of [Man11, Theorem 2.2.1]. The distance between the two hypersurfaces  $M_t = X(M^n, t)$  and  $N_t = Y(N^n, t)$  at time t, is given by

$$d(t) := \inf_{p \in M^n, q \in N^n} |X(p, t) - Y(q, t)|.$$

This function is locally Lipschitz in time, as the curvature is locally bounded and the two hypersurfaces move by mean curvature. Hence it is differentiable almost everywhere. Assume that t is a differentiability point. Since  $M^n$  is compact, d is actually a minimum. Suppose that d(t) > 0 and let  $(p_t, q_t) \in M^n \times N^n$  be points, where d(t) is attained. Differentiating |X(p,t) - Y(q,t)| with respect to  $v = v_1 \oplus v_2 \in T_{X(p,t)}M_t \bigoplus T_{Y(q,t)}N_t$  yields that

$$0 = \left\langle \frac{X(p_t, t) - Y(q_t, t)}{d(t)}, \nabla_{v_1} X(p_t, t) - \nabla_{v_2} Y(q_t, t) \right\rangle,$$

so that  $T_{X(p_t,t)}M_t$  and  $T_{Y(q_t,t)}N_t$  have to be parallel. Hence, we can write  $M_t$  and  $N_t$  locally around  $X(p_t,t)$  and  $Y(q_t,t)$  as graphs of two functions  $f,h:U\times(t-\varepsilon,t+\varepsilon)\to\mathbb{R}$ , where  $U\subset\mathbb{R}^n$ . After rotation, we can assume that  $\mathrm{span}(e_1,\ldots,e_n)\subset\mathbb{R}^{n+1}$  is such a tangent space with

$$X(p_t, t) = (0, f(0, t)), \quad Y(q_t, t) = (0, h(0, t)) \quad \text{and} \quad f(0, t) > h(0, t).$$

We calculate

$$\partial_t f = -H_M \langle \boldsymbol{\nu}_M, e_{n+1} \rangle = \Delta f - \frac{D_{ij} f D^i f D^j f}{1 + |D f|^2}$$

and

$$\partial_t h = -H_N \langle \boldsymbol{\nu}_N, e_{n+1} \rangle = \Delta h - \frac{D_{ij} h D^i h D^j h}{1 + |Dh|^2}.$$

The function f - h has a spatial minimum at x = 0 at time t. Hence,

$$\Delta f(0,t) - \Delta h(0,t) > 0$$
 and  $Df(0,t) = Dh(0,t) = 0$ 

and so

$$-\langle H_M(p_t, t) \nu_M(p_t, t) - H_N(q_t, t) \nu_N(q_t, t), e_{n+1} \rangle = \Delta f(0, t) - \Delta h(0, t) \ge 0.$$

Since

$$\frac{X(p_t, t) - Y(q_t, t)}{d(t)} = e_{n+1}$$

we obtain at  $(p_t, q_t)$  that

$$\begin{split} \partial_t |X(p,t) - Y(q,t)| \\ &= -\left\langle \frac{X(p_t,t) - Y(q_t,t)}{d(t)}, H_M(p_t,t) \boldsymbol{\nu}_M(p_t,t) - H_N(q_t,t) \boldsymbol{\nu}_N(q_t,t) \right\rangle \\ &= -\left\langle e_{n+1}, H_M(p_t,t) \boldsymbol{\nu}_M(p_t,t) - H_N(q_t,t) \boldsymbol{\nu}_N(q_t,t) \right\rangle \geq 0 \,. \end{split}$$

This holds for every minimum so that  $\partial_t d > 0$ .

**Proposition 1.9** (Preservation of embeddedness). If  $M_0$  is compact and embedded, then  $M_t$  is embedded for all  $t \in (0,T)$ .

In particular, let

$$m(t) := \max_{(p,s) \in M^n \times [0,t]} |A(p,s)|$$

and

$$l(p,q,t) := \int_{p}^{q} |\dot{\gamma}(s)|_{g(t)} ds$$
 for a minimizing geodesic  $\gamma$ 

and

$$\Omega_{\varepsilon}(t) := \{(p,q) \in M^n \times M^n \mid m(t)l(p,q,t) \le \varepsilon\}$$

for  $\varepsilon > 0$ . Then there exists  $\varepsilon > 0$  so that  $M_t$  is embedded on  $\Omega_{\varepsilon}(t)$  and

$$d(t) := \min_{(p,q) \in (M^n \times M^n) \backslash \Omega_{\varepsilon}(t)} d(p,q,t) \ge \min \left\{ d(0), \frac{\sin(\varepsilon)}{m(t)} \right\} \, .$$

*Proof.* We follow similar lines to [Man11, Proposition 2.2.7]. If the hypersurface  $M_0$  is embedded, then  $M_t$  is embedded for a small positive time, otherwise there is a sequence  $(p_i,q_i,t_i)_{i\in\mathbb{N}}$  with  $X(p_i,t_i)=X(q_i,t_i)$  and  $t_i\to 0$ . We have for a subsequence, that  $p_i\to p$  and  $q_i\to q$ . If  $p\neq q$  then X(p,0)=X(q,0), which is a contradiction. If p=q, by the smooth existence of the flow, there exists an open neighbourhood  $U\subset M^n$  of p so that the map  $X(\cdot,t)|_U$  is one-to-one for  $t\in [0,\varepsilon)$ , which is in contradiction. Define the monotone nondecreasing function

$$m(t) := \max_{(p,s) \in M^n \times [0,t]} |A(p,s)|$$

and we choose a smooth, monotone nondecreasing function  $m^*:[0,T)\to\mathbb{R}_+$  such that

$$m(t) \le m^*(t) \le 2m(t)$$

for every  $t \in [0,T)$ . Furthermore, define the geodesic intrinsic distance in the Riemannian manifold  $(M^n, g(t))$ 

$$l(p,q,t) := \int_{p}^{q} |\dot{\gamma}(s)|_{g(t)} ds$$
 for a minimizing geodesic  $\gamma$ 

and the extrinsic distances

$$d(p, q, t) := |X(p, t) - X(q, t)|$$
.

Consider the following inscribed and outscribed balls

$$B_{\text{in}}(p,t) := B_{1/m^*(t)} \left( X(p,t) - \frac{\nu(p,t)}{m^*(t)} \right)$$

and

$$B_{\text{out}}(p,t) := B_{1/m^*(t)} \left( X(p,t) + \frac{\boldsymbol{\nu}(p,t)}{m^*(t)} \right)$$

and the geodesic neighbourhood

$$U_{\varepsilon}(p,t) := \{ q \in M^n \mid m^*(t)l(p,q,t) \le \varepsilon \}.$$

Then there exists  $\varepsilon \in (0, \pi/2)$  so that

$$X(U_{\varepsilon}(p,t),t) \cap B_{\mathrm{in}}(p,t) = X(U_{\varepsilon}(p,t),t) \cap B_{\mathrm{out}}(p,t) = \emptyset$$

Consider the open set

$$\Omega_{\varepsilon}(t) := \{ (p, q) \in M^n \times M^n \mid m^*(t)l(p, q, t) \le \varepsilon \}$$

and the closed set

$$S(t) := \{ (p,q) \in M^n \times M^n \mid p \neq q \text{ and } X(p,t) = X(q,t) \}$$
.

For embedded  $M_t$ ,

$$\Omega_{\varepsilon}(t) \cap S(t) = \emptyset$$

and

$$d_{\partial\Omega_\varepsilon}(t):=\min_{(p,q)\in\partial\Omega_\varepsilon(t)}d(p,q,t)\geq\frac{2\sin(\varepsilon)}{m^*(t)}\,.$$

Assume that  $t_0 \in (0, T)$  is the first time where the flow is no more embedded. Since  $\Omega \cap S = \emptyset$  and  $\partial \Omega_{\varepsilon}(t_0)$  is compact,

$$\min_{t \in [0,t_0]} d_{\partial \Omega_\varepsilon}(t) = \frac{2\sin(\varepsilon)}{m^*(t_0)} \geq \frac{\sin(\varepsilon)}{m^*(t)} =: c > 0 \,.$$

Furthermore, set

$$d(t) := \min_{(p,q) \in (M^n \times M^n) \setminus \Omega_{\varepsilon}(t)} d(p,q,t).$$

Assume that there exists a time  $t_1 \in (0, t_0)$  so that  $d(t_1) < \min\{d(0), c\}$  for the first time. Then  $d(t_1)$  is attained at points  $(p_1, q_1) \in (M^n \times M^n) \setminus \Omega$ . A geometric argument analogous to the one of the comparison principle, Theorem 1.8, shows that  $\partial_t d(t) \geq 0$ . Hence

$$d(t) \ge \min\{d(0), c\} > 0$$

on  $[0, t_0]$ , which is a contradiction.

**Theorem 1.10** (Huisken, [Hui84, Corollary 3.6(ii)]). Let  $(M_t)_{t \in [0,T)}$  be a family of closed hypersurfaces moving by (MCF). Assume  $M_0 = X_0(M)$  closed and mean convex, i.e.  $H \ge 0$ . Then H > 0 for all  $t \in (0,T)$ .

Proof. See [Sch17c, Theorem 2.1.2]. That  $H \geq 0$  for  $t \geq 0$  follows from the evolution equation of H and the parabolic maximum principle, Theorem C.3. Assume that  $H(p_0, t_0) = 0$  for some  $t_0 > 0$ . The strong maximum principle then implies that H = 0 for all (p, t) and  $0 \leq t \leq t_0$ . But this is impossible since any closed hypersurface in  $\mathbb{R}^{n+1}$  has points where  $\lambda_1 > 0$ .

## 2. Homothetically shrinking solutions

**Definition 2.1** (Homothetically shrinking solutions, Brakke [Bra78, Appendix C]). Let  $\lambda: [t_0, T] \to \mathbb{R}_+$  be smooth and decreasing,  $\lambda(t_0) = 1$  and  $\lambda(T) = 0$ . Let  $x_0 \in \mathbb{R}^{n+1}$ . A homothetically shrinking solution  $X: M^n \times [t_0, T) \to \mathbb{R}^{n+1}$  to (MCF) satisfies

$$M_t = \lambda(t)(M_0 - x_0) + x_0$$

for all  $t \in [t_0, T)$ . This describes solutions of (MCF) which move by scaling about  $x_0$ .

**Remark 2.2.** See [Eck04, Examples 2.3(4)]. We can make the separation of variables ansatz

$$\tilde{X}(q,t) = \lambda(t)\tilde{X}(q,t_0)$$

for a family of embeddings  $\tilde{X}: M^n \times [t_0, T) \to \mathbb{R}^{n+1}$  with  $M_t = \tilde{X}(M^n, t)$  satisfying the evolution equation

$$\left(\partial_t \tilde{X}(q,t)\right)^{\perp} = \left\langle \partial_t \tilde{X}(q,t), \boldsymbol{\nu}(q,t) \right\rangle = \tilde{\mathbf{H}}(q,t)$$

for  $q \in M^n$ . In Remark 1.2, we saw that there are tangential diffeomorphisms  $\phi_t : M^n \to M^n$ ,  $t \in [t_0, T)$ , with

$$\tilde{X}(q,t) = X(\phi_t^{-1}(q), t)$$

for  $q \in M^n$ , where the embeddings  $X(\cdot,t): M^n \to \mathbb{R}^{n+1}$  satisfy (MCF). This says that, up to tangential diffeomorphisms, the radial or homothetic motion of the hypersurfaces  $M_t$  (described by  $\tilde{X}$ ) is equivalent to their normal motion along the mean curvature vector (described by X). For the shrinking sphere solution these two agree, but for the shrinking cylinder they differ. Since the mean curvature of the embeddings scales with factor  $1/\lambda(t)$  we deduce

$$\partial_t \lambda(t) \left( \tilde{X}(q, t_0) \right)^{\perp} = \left( \partial_t \tilde{X}(q, t) \right)^{\perp} = \tilde{\mathbf{H}}(q, t) = \frac{1}{\lambda(t)} \tilde{\mathbf{H}}(q, t_0)$$

for  $q \in M^n$ . From this we infer that

$$\alpha \equiv 2\lambda(t)\partial_t\lambda(t) = \partial_t\lambda^2(t)$$

is independent of t. We therefore obtain under the assumption  $\lambda(t_0) = 1$  that

$$\lambda(t) = \sqrt{1 + \alpha(t - t_0)}$$

for all t satisfying  $t > t_0 - 1/\alpha$ . Hence

$$\mathbf{H}(p,t) = \alpha \frac{\langle X(p,t), \boldsymbol{\nu}(p,t) \rangle}{2\lambda^2(t)}$$

for  $(p,t) \in M^n \times (-\infty, T)$ , where  $T = t_0 - 1/\alpha$ . This describes expanding homothetic solutions about 0 for  $\alpha > 0$  and contracting homothetic solutions about 0 for  $\alpha < 0$ . Let us concentrate on  $\alpha < 0$ . If we set  $\lambda(T) = 0$  for  $T > t_0$ , which requires the hypersurface to disappear at time T, then  $\alpha = -1/(T - t_0)$  and thus

$$\lambda(t) = \sqrt{\frac{T - t}{T - t_0}}$$

and

$$\mathbf{H}(p,t) = \frac{\langle X(p,t), \boldsymbol{\nu}(p,t) \rangle}{2(T-t)}$$

for  $(p,t) \in M^n \times (-\infty,T)$ .

**Lemma 2.3.** Let  $(M_t)_{t \in (-\infty,0)}$  be an ancient solution of MCF. Then

$$H(x) = \frac{\langle x, \boldsymbol{\nu}(x) \rangle}{-2t}$$

for all  $x \in M_t$  and t < 0 if and only if  $M_t = \sqrt{-t}M_{-1}$  for all t < 0.

*Proof.* Let  $M_t = \sqrt{-t}M_{-1}$  for all t < 0. Then  $H(x) = \langle x, \boldsymbol{\nu}(x) \rangle / (-2t)$  for all  $x \in M_t$  and t < 0 follows by Remark 2.2.

On the other hand, let  $H(x) = \langle x, \boldsymbol{\nu}(x) \rangle / (-2t)$  for all  $x \in M_t$  and t < 0. Then

$$\langle \Delta_{M_t} X(p,t), \boldsymbol{\nu}(p,t) \rangle = -H(p,t) = -\frac{\langle X(p,t), \boldsymbol{\nu}(p,t) \rangle}{-2t}$$

and thus up to tangential motion  $X(p,t) = \sqrt{-2t}X(p,t_0)$ .

### 2.1. Hypersurfaces.

**Theorem 2.4** (Huisken, [Hui90, Theorem 4.1] and [Hui93]). Let  $M \subset \mathbb{R}^{n+1}$  be a smooth, complete, embedded, mean convex hypersurface such that  $H(x) = \langle x, \nu \rangle / 2$  at every  $x \in M$  and there exists a constant C > 0 such that  $|A| + |\nabla A| \leq C$  and  $\mu^n(M \cap B_R) \leq Ce^R$ , for every ball of radius R > 0 in  $\mathbb{R}^{n+1}$ . Then, up to a rotation in  $\mathbb{R}^{n+1}$ , M is of the form  $\mathbb{S}^m_{\sqrt{2m}} \times \mathbb{R}^{n-m}$  for  $m = 0, 1, \ldots, n$ .

*Proof.* See [Man11, Proposition 3.4.1]. We scale M by the factor 1/2 so that  $H(x) = \langle x, \boldsymbol{\nu}(x) \rangle$  at every  $x \in M$ . By covariant differentiation of the equation  $H = \langle x, \boldsymbol{\nu} \rangle$  in an orthonormal frame  $\{\boldsymbol{\tau}_1, \dots, \boldsymbol{\tau}_n\}$  on M we get by the Weingarten equations  $\nabla_i \boldsymbol{\nu} = \partial_i \boldsymbol{\nu} = h_i^j \partial_j x$  that

$$\nabla_j H = \langle x, \nabla_j \boldsymbol{\nu} \rangle = \langle x, \partial_k x \rangle h_j^k$$

and by the Gauss equations  $\nabla_i \nabla_j x = -h_{ij} \nu$  and Codazzi equations  $\nabla_k h_{ij} = \nabla_k h_{ji} = \nabla_j h_{ik}$  at one fixed point where the Christoffel symbols vanish, that

$$\nabla_{i}\nabla_{j}H = g_{ik}h_{j}^{k} + \langle x, \nabla_{i}\nabla_{k}x\rangle h_{j}^{k} + \langle x, \partial_{k}K\rangle \nabla_{i}h_{j}^{k}$$

$$= h_{ij} + \langle x, \nu\rangle h_{ik}h_{j}^{k} + \langle x, \partial_{k}x\rangle g^{kl}\nabla_{i}h_{jl}$$

$$= h_{ij} - Hh_{ik}h_{j}^{k} + \langle x, \partial_{k}x\rangle g^{kl}\nabla_{l}h_{ij}$$

$$= h_{ij} - Hh_{ik}h_{j}^{k} + \langle x, \nabla h_{ij}\rangle. \tag{1}$$

Contracting with  $g^{ij}$  we have

$$\Delta H = H \left( 1 - |A|^2 \right) + \langle x, \nabla H \rangle. \tag{2}$$

From equation (2) and the strong maximum principle for elliptic equations, Theorem C.1, we see that, since M satisfies  $H \geq 0$  by assumption and

$$\Delta H \le H + \langle x, \nabla H \rangle$$

we must either have that H=0 or H>0 on all M. Contracting (1) with  $h^{ij}$ , we have

$$h^{ij}\nabla_i\nabla_j H = |A|^2 - H\operatorname{tr}(A^3) + \frac{\langle x, \nabla |A|^2\rangle}{2},$$

which implies, by Simons' identity

$$\Delta h_{ij} = \nabla_i \nabla_j H + H h_{ik} h_j^k - |A|^2 h_{ij}$$

that

$$\begin{split} \Delta |A|^2 &= \Delta (h^{ij}h_{ij}) = h^{ij}\Delta h_{ij} + 2g^{mn}\nabla_m h^{ij}\nabla_n h_{ij} + h_{ij}\Delta h^{ij} \\ &= h^{ij}\Delta h_{ij} + 2g^{mn}g^{ki}g^{lj}\nabla_m h_{kl}\nabla_n h_{ij} + h_{ij}g^{ki}g^{jl}\Delta h_{kl} \\ &= 2h^{ij}\left(\nabla_i\nabla_j H + H h_{ik}h_j^k - |A|^2 h_{ij}\right) + 2g^{mn}\nabla_m h_i^i\nabla_n h_i^l \\ &= 2|A|^2 - 2H\operatorname{tr}(A^3) + \langle x, \nabla |A|^2 \rangle + 2H\operatorname{tr}(A^3) - 2|A|^4 + 2|\nabla A|^2 \\ &= 2|A|^2(1 - |A|^2) + \langle x, \nabla |A|^2 \rangle + 2|\nabla A|^2 \,. \end{split}$$

Assume that H=0. As M is complete and x is a tangent vector field on M by the equation  $\langle x, \boldsymbol{\nu} \rangle = 0$ , for every point  $x \in M$  there is a unique solution of the ODE

$$\gamma'(s) = x(\gamma(s)) = \gamma(s)$$

passing through x and contained in M for every  $s \in \mathbb{R}$ , but such solution is simply the line in  $\mathbb{R}^{n+1}$  passing through x and the origin. Thus, M has to be a cone and being smooth the only possibility is a hyperplane through the origin of  $\mathbb{R}^{n+1}$ .

Assume that H > 0 everywhere (so dividing by H and |A| is allowed). For R > 0, define

$$\eta_R = \boldsymbol{\nu}_{\partial(M \cap B_R(0))}$$

to be the outward unit conormal to  $M \cap B_R(0)$  along  $\partial(M \cap B_R(0))$ , which is a smooth boundary for almost every R > 0 (by Sard's theorem, see homework or Corollary B.3). Then, supposing that R belongs to the set  $\mathcal{R} \subset \mathbb{R}^+$  of the regular values of the function  $|\cdot|$  restricted to  $M \subset \mathbb{R}^{n+1}$ , from equation (2) and the divergence theorem, Theorem A.2, we compute

$$\begin{split} \varepsilon_R &= \int_{\partial(M\cap B_R(0))} |A| \langle \nabla H, \eta_R \rangle \exp\left(-\frac{R^2}{2}\right) \, d\mu^{n-1} \\ &= \int_{M\cap B_R(0)} |A| \Delta H \exp\left(-\frac{|x|^2}{2}\right) + \left\langle \nabla \left(|A| \exp\left(-\frac{|x|^2}{2}\right)\right), \nabla H \right\rangle \, d\mu^n \\ &= \int_{M\cap B_R(0)} \left(|A| H \left(1 - |A|^2\right) + |A| \langle x, \nabla H \rangle\right) \exp\left(-\frac{|x|^2}{2}\right) \, d\mu^n \\ &+ \int_{M\cap B_R(0)} \left(\frac{1}{2|A|} \langle \nabla |A|^2, \nabla H \rangle - |A| \langle x, \nabla H \rangle\right) \exp\left(-\frac{|x|^2}{2}\right) \, d\mu^n \\ &= \int_{M\cap B_R(0)} \left(|A| H \left(1 - |A|^2\right) + \frac{1}{2|A|} \langle \nabla |A|^2, \nabla H \rangle\right) \exp\left(-\frac{|x|^2}{2}\right) \, d\mu^n \end{split}$$

and similarly

$$\begin{split} \delta_R &= \int_{\partial(M\cap B_R(0))} \frac{H}{|A|} \langle \nabla |A|^2, \eta_R \rangle \exp\left(-\frac{R^2}{2}\right) \, d\mu^{n-1} \\ &= \int_{M\cap B_R(0)} \frac{H}{|A|} \Delta |A|^2 \exp\left(-\frac{|x|^2}{2}\right) + \left\langle \nabla \left(\frac{H}{|A|} \exp\left(-\frac{|x|^2}{2}\right)\right), \nabla |A|^2 \right\rangle \, d\mu^n \\ &= \int_{M\cap B_R(0)} \left(2|A|H \left(1-|A|^2\right) + \frac{2H|\nabla A|^2}{|A|} \right. \\ &\qquad \qquad \left. + \frac{H}{|A|} \langle x, \nabla |A|^2 \rangle \right) \exp\left(-\frac{|x|^2}{2}\right) \, d\mu^n \\ &\qquad \qquad + \int_{M\cap B_R(0)} \left(\frac{\langle \nabla H, \nabla |A|^2 \rangle}{|A|} - \frac{H|\nabla |A|^2|^2}{2|A|^3} - \frac{H}{|A|} \langle x, \nabla |A|^2 \rangle \right) \exp\left(-\frac{|x|^2}{2}\right) \, d\mu^n \\ &= \int_{M\cap B_R(0)} \left(2|A|H \left(1-|A|^2\right) + \frac{2H|\nabla A|^2}{|A|} + \frac{\langle \nabla H, \nabla |A|^2 \rangle}{|A|} \right. \\ &\qquad \qquad - \frac{H|\nabla |A|^2|^2}{2|A|^3} \right) \exp\left(-\frac{|x|^2}{2}\right) \, d\mu^n \; . \end{split}$$

Hence,

$$\begin{split} \sigma_R &= 2\delta_R - 4\varepsilon_R \\ &= \int_{M \cap B_R(0)} \left( \frac{4H|\nabla A|^2}{|A|} - \frac{H|\nabla |A|^2|^2}{|A|^3} \right) \exp\left( -\frac{|x|^2}{2} \right) d\mu^n \\ &= \int_{M \cap B_R(0)} \left( 4|A|^2 |\nabla A|^2 - |\nabla |A|^2 |^2 \right) \frac{H}{|A|^3} \exp\left( -\frac{|x|^2}{2} \right) d\mu^n \,. \end{split}$$

As we have

$$4|A|^2|\nabla A|^2 \ge |\nabla |A|^2|^2$$

the quantity  $\sigma_R$  is nonnegative and nondecreasing in R. If now we show that

$$\liminf_{R \to \infty} \sigma_R = 0$$

we can conclude that, at every point of M,

$$4|A|^2|\nabla A|^2 = |\nabla |A|^2|^2. \tag{3}$$

Getting back to the definitions of  $\varepsilon_R$  and  $\delta_R$ , we have

$$|\sigma_R| = \left| -2 \int_{\partial(M \cap B_R(0))} \frac{H}{|A|} \langle \nabla |A|^2, \eta \rangle \exp\left(-\frac{R^2}{2}\right) d\mu^{n-1} \right|$$

$$+4 \int_{\partial(M \cap B_R(0))} |A| \langle \nabla H, \eta \rangle \exp\left(-\frac{R^2}{2}\right) d\mu^{n-1}$$

$$\leq 4 \exp\left(-\frac{R^2}{2}\right) \int_{\partial(M \cap B_R(0))} \left(\frac{H}{|A|} |\nabla |A|^2 | + |A| |\nabla H|\right) d\mu^{n-1}$$

$$\leq 8 \exp\left(-\frac{R^2}{2}\right) \int_{\partial(M \cap B_R(0))} (H|\nabla A| + |A| |\nabla H|) d\mu^{n-1}$$

$$\leq C \exp\left(-\frac{R^2}{2}\right) \mu^{n-1} (\partial(M \cap B_R(0))),$$

by the estimates on A and  $\nabla A$  in the hypotheses. Assume that the lefthand side does not go to zero. That is, suppose that for every R belonging to the set  $\mathcal{R} \subset \mathbb{R}^+$  (which is of full measure) and R larger than some  $R_0 > 0$  we have

$$\mu^{n-1}(\partial(M \cap B_R(0))) \ge \delta \exp\left(\frac{R^2}{2}\right) \ge \delta R \exp\left(\frac{R^2}{4}\right)$$

for some constant  $\delta > 0$ . Recall the area formula and divergence theorem, Theorems A.1 and A.2. As the function

$$R \mapsto \mu^n(M \cap B_R(0))$$

is monotone and continuous from the left and actually continuous at every value  $R \in \mathcal{R}$ , we can differentiate it almost everywhere in  $\mathbb{R}^+$  and we have, for  $R_0 < r < R$ .

$$\mu^{n}(M \cap B_{R}(0)) - \mu^{n}(M \cap B_{r}(0)) = \int_{r}^{R} \frac{d}{d\xi} \mu^{n}(M \cap B_{\xi}(0)) d\xi$$

$$= \int_{r}^{R} \int_{M \cap B_{\xi}(0)} \operatorname{div}_{M \cap B_{\xi}(0)} \eta_{\xi} d\mu^{n-1} d\xi$$

$$= -\int_{r}^{R} \int_{M \cap B_{\xi}(0)} \langle \eta_{\xi}, \mathbf{H}_{M \cap B_{\xi}(0)} \rangle d\mu^{n-1} d\xi$$

$$+ \int_{r}^{R} \int_{\partial (M \cap B_{\xi}(0))} \langle \eta_{\xi}, \eta_{\xi} \rangle d\mu^{n-1} d\xi$$

$$= \int_{r}^{R} \int_{\partial (M \cap B_{\xi}(0))} d\mu^{n-1} d\xi$$

$$\geq \delta \int_{r}^{R} \xi \exp\left(\frac{\xi^{2}}{4}\right) d\xi = 2\delta \left(\exp\left(\frac{R^{2}}{4}\right) - \exp\left(\frac{r^{2}}{4}\right)\right).$$

Then

$$\mu^n(M \cap B_R(0))e^{-R} \to \infty$$
,

for  $R \to \infty$ , in contradiction with the hypotheses of the theorem. Hence, the

$$\liminf_{R\to\infty,R\in\mathcal{R}} \exp\biggl(-\frac{R^2}{2}\biggr)\,\mu^{n-1}(\partial(M\cap B_R(0))) = 0\,.$$

It follows that the same holds for  $|\sigma_R|$  and equation (3) is proved. By Cauchy–Schwarz,

$$4|A|^2|\nabla A|^2 = |\nabla |A|^2|^2 = 4|A\nabla A|^2 \le 4|A|^2|\nabla A|^2$$

or in coordinates

$$\begin{aligned} 4h_j^i h_i^j \nabla_k h_n^m \nabla^k h_m^n &= \nabla_k (h_j^i h_i^j) \nabla^k (h_n^m h_m^n) \\ &= 4h_j^i h_n^m \nabla_k h_i^j \nabla^k h_m^n \leq 4h_j^i h_i^j \nabla_k h_n^m \nabla^k h_m^n \end{aligned}$$

with equality if and only if A and  $\nabla A$  are linearly dependent, that is, at every point there exist constants  $c_k$  such that

$$\nabla_k h_{ij} = c_k h_{ij}$$

for every i, j. Contracting this equation with the metric  $g^{ij}$  and with  $h^{ij}$  we get

$$\nabla_k H = c_k H$$
 and  $\nabla_k |A|^2 = 2c_k |A|^2$ ,

hence

$$\nabla_k \log H = c_k$$
 and  $\nabla_k \log |A|^2 = 2c_k$ .

This implies

$$\nabla_k \log \left( \frac{H}{|A|} \right) = 0$$
 so that  $|A| = \alpha H$ 

for some constant  $\alpha > 0$ . By connectedness this relation has to hold globally on M. Suppose now that at a point  $|\nabla H| \neq 0$ , then

$$\nabla_k h_{ij} = c_k h_{ij} = \frac{\nabla_k H}{H} h_{ij} \tag{4}$$

which is a symmetric 3-tensor by the Codazzi equations, hence

$$h_{ij}\nabla_k H = h_{ik}\nabla_i H$$

at one point, where the Christoffel symbols vanish. Computing then in normal coordinates with an orthonormal basis  $\{\boldsymbol{\tau}_1,\ldots,\boldsymbol{\tau}_n\}$  such that  $\boldsymbol{\tau}_1=\nabla H/|\nabla H|$ , we have with  $g^{ij}=\delta^{ij}$ ,

$$0 = |h_{ij}\nabla_{k}H - h_{ik}\nabla_{j}H|^{2}$$

$$= (h_{ij}\nabla_{k}H - h_{ik}\nabla_{j}H)g^{il}g^{jm}g^{kn}(h_{lm}\nabla_{n}H - h_{ln}\nabla_{m}H)$$

$$= 2|\nabla H|^{2}|A|^{2} - 2g^{il}g^{jm}g^{kn}h_{ij}h_{ln}\nabla_{k}H\nabla_{m}H$$

$$= 2|\nabla H|^{2}|A|^{2} - 2g^{il}h_{i}^{m}h_{i}^{k}\nabla_{k}H\nabla_{m}H$$

$$= 2|\nabla H|^{2}|A|^{2} - 2g^{il}h_{i}^{1}h_{i}^{1}\nabla_{1}H\nabla_{1}H$$

$$= 2|\nabla H|^{2}|A|^{2} - 2g^{il}h_{i}^{1}h_{i}^{1}\nabla_{1}H\nabla_{1}H$$

$$= 2|\nabla H|^{2}\left(|A|^{2} - \sum_{i=1}^{n}(h_{i}^{1})^{2}\right).$$

Hence,  $|A|^2 = \sum_{i=1}^n (h_i^1)^2$  and

$$|A|^2 = (h_1^1)^2 + 2\sum_{i=2}^n (h_i^1)^2 + \sum_{i,j\neq 1}^n (h_i^j)^2$$

so  $h_j^i = 0$  unless i = j = 1, which means that A has rank one. Thus, we have two possible (not mutually excluding) situations at every point of M, either A has rank one or  $\nabla H = 0$ .

If ker  $A \equiv \emptyset$  on M, A must have rank at least two as we assumed  $n \geq 2$ , then we have  $\nabla H = 0$  which implies  $\nabla A = 0$  and

$$h_{ij} = H h_{ik} h_j^k = H h_{ik} g^{kl} h_{lj}$$

by equation (1). This means that for an eigenvalue  $\lambda_m$  with eigenvector  $\xi_m$ ,

$$h_{ij}\xi_m^j = Hh_{ik}g^{kl}h_{lj}\xi_m^j = Hh_{ik}g^{kl}\lambda_m g_{lj}\xi_m^j = \lambda_m Hh_{ij}\xi_m^j$$

so that all the eigenvalues of A are 0 or 1/H. As the kernel is empty

$$H = \sum_{i=1}^{n} \lambda_m = \frac{n}{H}$$

so that

$$H = \sqrt{n}$$
 and  $h_{ij} = \frac{g_{ij}}{\sqrt{n}}$ .

Then, the complete hypersurface M has to be the sphere  $\mathbb{S}_{\sqrt{n}}^n$ , indeed we compute

$$\Delta |x|^2 = \Delta |x|^2 = 2\nabla \langle x, \nabla x \rangle = 2n + 2\langle x, \Delta x \rangle$$
$$= 2n - 2H\langle x, \boldsymbol{\nu} \rangle = 2n - 2H^2 = 0,$$

by means of the structural equation  $H = \langle x, \boldsymbol{\nu} \rangle$ . Hence,  $|x|^2$  is a harmonic function on M. Looking at the point of M of minimum distance from the origin, by the strong maximum principle for elliptic equations, Theorem C.1, it must be constant on M and  $M = \mathbb{S}^n_{\sqrt{n}}$ .

Let now  $\ker A(x) \neq \emptyset$  at some point  $x \in M$ , with  $\dim \ker A(x) = (n-m)$  and 0 < m < n (as A is never zero), and let

$$v_1(x), \dots, v_{n-m}(x) \in T_x M \subset \mathbb{R}^{n+1}$$

be a family of unit orthonormal tangent vectors spanning ker A(x), that is,

$$h_{ij}(x)v_k^j(x) = 0$$

for k = 1, ..., n - m. By (4), the geodesic  $\gamma(s)$  from  $x \in M$  (M is complete) with initial velocity  $\partial_s \gamma(0) = v_k(x)$  satisfies

$$\nabla_{\partial_s \gamma} (h_{ij} \partial_s \gamma^j) = \frac{\langle \nabla H, \partial_s \gamma \rangle}{H} h_{ij} \partial_s \gamma^j$$

hence, by Gronwalls lemma there holds

$$h_{ij}(\gamma(s))\partial_s \gamma^j(s) = h_{ij}(\gamma(0))\partial_s \gamma^j(0) \exp\left(\int_0^s \frac{\langle \nabla H, \partial_\sigma \gamma \rangle}{H} d\sigma\right) = 0$$

for every  $s \in \mathbb{R}$ . Since  $\gamma$  is a geodesic in M,  $\partial_s^2 \gamma(s) \in (T_{\gamma(s)}M)^{\perp}$ , that is, the normal to the curve in  $\mathbb{R}^{n+1}$  is also the normal to M, then letting  $\kappa$  be the curvature of  $\gamma$  in  $\mathbb{R}^{n+1}$ , we have

$$\kappa = -\langle \boldsymbol{\nu}_{M}, \partial_{s}^{2} \gamma \rangle = h_{ij} \partial_{s} \gamma^{i} \partial_{s} \gamma^{j} = 0 \,,$$

thus  $\gamma$  is a straight line in  $\mathbb{R}^{n+1}$  and

$$x + \ker A(x) \subset M$$
,

where  $x + \ker A(x) \subset \mathbb{R}^{n+1}$  is an (n-m)-dimensional affine subspace. Let now  $\sigma(s)$  be a geodesic from x to another point y parametrized by arclength and extend by parallel transport the vectors  $v_k(x)$ ,  $k = 1, \ldots, n-m$ , along  $\sigma$ , then

$$\nabla_{\partial_s \sigma} (h_{ij} v_k^j) = \frac{\langle \nabla H, \partial_s \sigma \rangle}{H} h_{ij} v_k^j$$

and again by Gronwalls lemma it follows that  $h_{ij}(\gamma(s))v_k^j(\gamma(s))=0$  for every  $s \in \mathbb{R}$  and  $k=1,\ldots,n-m$ , in particular  $v_k(y) \in \ker A(y)$ . Hence,

$$\dim \ker A \equiv n - m$$

on M with 0 < m < n (as A is never zero) and all the affine (n-m)-dimensional subspaces  $x + \ker A(x) \subset \mathbb{R}^{n+1}$  are contained in M for every  $x \in M$ , that is,

$$M + \ker(M) \subset M$$
.

Moreover, as  $h_{ij}v_k^j=0$  along the geodesic  $\sigma$ , we have

$$D_{\partial_s \sigma}^{\mathbb{R}^{n+1}} v_k = \nabla_{\partial_s \sigma} v_k + \langle \nabla_{\partial_s \sigma} v_k, \boldsymbol{\nu}_M \rangle \boldsymbol{\nu}_M = -h_{ij} v_k^j \partial_s \sigma^i \boldsymbol{\nu}_M = 0,$$

so the extended vectors  $v_k$  are constant in  $\mathbb{R}^{n+1}$ , which means that the parallel extension is independent of the geodesic  $\sigma$ , that the subspaces  $\ker A(x)$  are all a common (n-m)-dimensional vector subspace of  $\mathbb{R}^{n+1}$  and

$$M = M + \ker A$$
.

Let  $x \in M$ . Then there exists  $y \in M \cap (\ker A)^{\perp}$  and  $v \in \ker A$  so that

$$x = y + v$$
.

Define  $f: M \to \ker A$  by

$$f(x) = v$$
.

By Sards theorem, Corollary B.3, there exists a vector  $v \in \ker A$  such that

$$N(v) := f^{-1}(v) = M \cap \left(v + (\ker A)^{\perp}\right)$$

is a smooth, complete m-dimensional submanifold of  $\mathbb{R}^{n+1}$ . Since  $M=M+\ker A$ , N(v)=N(w) for all  $v,w\in\ker A$  and

$$M = N \times \ker A$$
.

This implies that

$$L := N(0) = M \cap (\ker A)^{\perp}$$

is a smooth, complete m-dimensional submanifold of  $(\ker A)^{\perp} = \mathbb{R}^{m+1}$  with

$$M = L \times \ker A$$
.

Moreover, as ker A is in the tangent space to every point of L, the normal  $\nu_M$  to M at a point of L stays in  $(\ker A)^{\perp}$  so it must coincide with the normal  $\nu_L$  to L in  $(\ker A)^{\perp}$ , then a simple computation shows that the mean curvature  $H_M$  of M at the points of L is equal to the mean curvature  $H_L$  of L as a hypersurface of  $(\ker A)^{\perp} = \mathbb{R}^{m+1}$ . This shows that L is a hypersurface in  $\mathbb{R}^{m+1}$  satisfying  $H_L(z) = \langle z, \nu_L(z) \rangle$  for every  $z \in L$ . Finally, as by construction the second fundamental form of L has empty kernel, by the previous discussion we have  $L = \mathbb{S}^m_{\sqrt{m}}$  and  $M = \mathbb{S}^m_{\sqrt{m}} \times \mathbb{R}^{n-m}$  which proves the claim.

**Theorem 2.5** (Colding–Minicozzi, [CM12, Theorem 10.1]). If  $M^n$ , for  $n \geq 2$ , is an embedded hypersurface in  $\mathbb{R}^{n+1}$ , with non-negative mean curvature, satisfying  $H = \langle x, \nu \rangle / 2$ , then  $M^n$  is of the form  $\mathbb{S}^m_{\sqrt{2m}} \times \mathbb{R}^{n-m}$  for  $m = 0, 1, \ldots, n$ .

## 2.2. Curves.

**Theorem 2.6** (Abresch-Langer, [AL86]). Let  $\Sigma \subset \mathbb{R}^2$  be a smooth, complete, embedded curve satisfying  $\kappa(x) = \langle x, \boldsymbol{\nu}(x) \rangle / 2$  for every  $x \in \Sigma$ . Then  $\Sigma$  is either the line through the origin or the  $\mathbb{S}^1_{\sqrt{2}}$ .

*Proof.* See [Man11, Proposition 3.4.1]. We scale the curve by the factor 1/2 so that  $\kappa = \langle x, \boldsymbol{\nu} \rangle$  for every  $x \in \Sigma$ . Fixing a reference point on a curve  $\Sigma = X(I)$ ,  $I \in \{\mathbb{S}^1, \mathbb{R}\}$ , we have an arclength parameter s which gives a unit tangent vector field  $\boldsymbol{\tau} = \partial_s X$  and a unit normal vector field  $\boldsymbol{\nu} = (\boldsymbol{\tau}_2, -\boldsymbol{\tau}_1)$ , which is the clockwise rotation of  $\pi/2$  in  $\mathbb{R}^2$  of the vector  $\boldsymbol{\tau}$ . Then the curvature is given by

$$\kappa = -\langle \partial_s \boldsymbol{\tau}, \boldsymbol{\nu} \rangle = \langle \boldsymbol{\tau}, \partial_s \boldsymbol{\nu} \rangle$$

so that

$$\partial_s \boldsymbol{\nu} = \kappa \boldsymbol{\tau}$$
 and  $\partial_s \boldsymbol{\tau} = -\kappa \boldsymbol{\nu}$ .

The relation  $\kappa = \langle x, \boldsymbol{\nu} \rangle$  implies the ODE for the curvature

$$\partial_s \kappa = \langle \boldsymbol{\tau}, \boldsymbol{\nu} \rangle + \langle x, \partial_s \boldsymbol{\nu} \rangle = \kappa \langle x, \boldsymbol{\tau} \rangle.$$

Suppose that at some point  $\kappa = 0$ , then also  $\partial_s \kappa = 0$  at the same point. Hence, by the uniqueness theorem for ODEs we conclude that  $\kappa$  is identically zero so that  $\Sigma$  is a line. Since  $\langle x, \boldsymbol{\nu} \rangle = 0$  for every  $x \in \Sigma$ , we conclude that  $0 \in \Sigma$ . So we suppose that  $\kappa$  is always nonzero and possibly reversing the orientation of the curve, we

assume that  $\kappa > 0$  at every point, that is, the curve is strictly convex. Computing the derivative of  $|X|^2$ ,

$$\partial_s |X|^2 = 2\langle X, \boldsymbol{\tau} \rangle = 2 \frac{\partial_s \kappa}{\kappa} = 2 \partial_s \log \kappa$$

we get

$$\kappa = C \exp\left(\frac{|x|^2}{2}\right)$$

for some constant C > 0. Hence,  $\kappa$  is bounded from below by C > 0. Since  $\Sigma$  is convex, we can consider the coordinate  $\vartheta = \arccos\langle e_1, \nu \rangle$ . (Note that  $\vartheta$  is only locally continuous and jumps after a complete round). We have  $\partial_s \vartheta = \kappa$  as well as

$$\partial_{\vartheta}\kappa = \frac{\partial_{s}\kappa}{\kappa} = \langle x, \boldsymbol{\tau} \rangle$$
 and  $\partial_{\vartheta}^{2}\kappa = \frac{\partial_{s}\partial_{\vartheta}\kappa}{\kappa} = \frac{1 - \kappa\langle x, \boldsymbol{\nu} \rangle}{\kappa} = \frac{1}{\kappa} - \kappa$ . (5)

Multiplying both sides of the last equation by  $2\partial_{\vartheta}\kappa$  we get

$$0 = 2\partial_{\vartheta}\kappa\partial_{\vartheta}^{2}\kappa + 2\kappa\partial_{\vartheta}\kappa - \frac{2\partial_{\vartheta}\kappa}{\kappa} = \partial_{\vartheta}\left((\partial_{\vartheta}\kappa)^{2} + \kappa^{2} - \log\kappa^{2}\right),$$

so that.

$$(\partial_{\vartheta}\kappa)^2 + \kappa^2 - \log \kappa^2 \equiv E \ge 1$$

along all the curve. We have E=1 if and only if  $\kappa^2\equiv 1$  along the curve, which is the unit circle centered at the origin of  $\mathbb{R}^2$ . When E>1, it follows that  $\kappa$  is uniformly bounded from above, hence recalling that  $\kappa=C\exp(|x|^2/2)$ ,

$$\Sigma \subset B_R(0)$$

for some R > 0 and by the embeddedness and completeness hypotheses,  $\Sigma$  must be closed, simple and strictly convex, as  $\kappa > 0$  at every point.

Suppose that  $\Sigma$  is not a line. We follow the lines of [GH86, Lemma 5.7.9] and [Pih98, Lemma 7.23]. The system

$$\left\{1, \sqrt{2}\cos(n\vartheta), \sqrt{2}\sin(n\vartheta)\right\}_{n\in\mathbb{Z}} \tag{6}$$

forms an orthonormal basis of the periodic functions in the Hilbert space  $C^2([0, 2\pi])$  with respect to the  $L^2$ -inner product (see e.g. [HL99, p. 124]). We have  $ds_t = d\vartheta/\kappa$  so that

$$\int_{\mathbb{S}^1} \frac{\sin(\vartheta)}{\kappa} d\vartheta = \int_{\mathbb{S}^1_{R_t}} \sin\left(\frac{s}{R_t}\right) ds_t = \cos(2\pi) - \cos(0) = 1 - 1 = 0$$

and

$$\int_{\mathbb{S}^1} \frac{\cos(\vartheta)}{\kappa} d\vartheta = \int_{\mathbb{S}^1_{R_t}} \cos\left(\frac{s}{R_t}\right) ds_t = \sin(2\pi) - \sin(0) = 0.$$

Furthermore, integration by parts yields

$$\begin{split} 0 &= \int_{\mathbb{S}^1} \frac{\sin(\vartheta)}{\kappa} \, d\vartheta \int_{\mathbb{S}^1} \frac{1}{\kappa} \frac{\partial \cos}{\partial \vartheta} (\vartheta) \, d\vartheta \\ &= -\int_{\mathbb{S}^1} \cos(\vartheta) \frac{\partial}{\partial \vartheta} \bigg( \frac{1}{\kappa} \bigg) \, d\vartheta = \int_{\mathbb{S}^1} \cos(\vartheta) \frac{1}{\kappa^2} \frac{\partial \kappa}{\partial \vartheta} \, d\vartheta \end{split}$$

and

$$\begin{split} 0 &= -\int_{\mathbb{S}^1} \frac{\cos(\vartheta)}{\kappa} \, d\vartheta = \int_{\mathbb{S}^1} \frac{1}{\kappa} \frac{\partial \sin}{\partial \vartheta} (\vartheta) \, d\vartheta \\ &= -\int_{\mathbb{S}^1} \sin(\vartheta) \frac{\partial}{\partial \vartheta} \bigg( \frac{1}{\kappa} \bigg) \, d\vartheta = \int_{\mathbb{S}^1} \sin(\vartheta) \frac{1}{\kappa^2} \frac{\partial \kappa}{\partial \vartheta} \, d\vartheta \, . \end{split}$$

Additionally, we have

$$0 = -\int_{\mathbb{S}^1} \frac{\partial}{\partial \vartheta} \left( \frac{1}{\kappa} \right) d\vartheta = \int_{\mathbb{S}^1} \frac{1}{\kappa^2} \frac{\partial \kappa}{\partial \vartheta} d\vartheta \,.$$

Hence,  $1/\kappa^2 \frac{\partial}{\partial \vartheta} \kappa$  is orthogonal to the first five basis functions of the basis (6). Since all the other basis functions are zero at at least four points in  $[0, 2\pi]$  with distance  $\leq \pi/2$ , there exists a number  $i_0 \geq 4$  and points  $\vartheta_i \in \mathbb{S}^1$ ,  $i \in \{0, \dots, i_0\}$ , so that

$$\left(\frac{1}{\kappa^2}\frac{\partial \kappa}{\partial \vartheta}\right)(\vartheta_i, \tau) = 0$$

and

$$|\vartheta_i - \vartheta_{i+1}| \le \frac{\pi}{2}$$

for  $i \in \{0, ..., i_0 - 1\}$  and

$$|\vartheta_{i_0} - (2\pi + \vartheta_0)| \le \frac{\pi}{2}.$$

Since  $1/\kappa^2 \frac{\partial}{\partial \theta} \kappa$  is periodic on  $[0, 2\pi]$ ,  $i_0$  is odd. Define the intervals

$$I_i := [\vartheta_i, \vartheta_{i+1}]$$

for  $i \in \{0, ..., i_0 - 1\}$  and

$$I_{i_0} := [0, \vartheta_0] \cup [\vartheta_{i_0}, 2\pi).$$

Then  $|I_i| \leq \pi/2$  for all  $i \in \{1, \ldots, i_0\}$ . Since  $\partial_{\vartheta}^2 \kappa = 1/\kappa - \kappa$ , it holds that  $\partial_{\vartheta}^2 \kappa \neq 0$  when  $\partial_{\vartheta} \kappa = 0$ , otherwise this second-order ODE for  $\kappa$  would imply  $\partial_{\vartheta} \kappa = 0$  everywhere, hence  $\kappa = 1$  identically and we would be in the case of the unit circle. Suppose that  $\Sigma$  is neither a line nor a circle. By looking at the equation for the curvature (5) we can see easily that  $\kappa < 1$  at a local minimum and  $\kappa > 1$  at a local maximum. Suppose now that  $\kappa(0)$  is a local maximum and  $\kappa(\vartheta_0)$  is the first subsequent critical value for  $\kappa$  for  $\vartheta_0 \leq \pi/2$  by the above. Then the curvature is strictly decreasing in the interval  $[0, \vartheta_0]$ . Also  $\kappa(\vartheta_0) < 1$  must be a local minimum of the curvature, as every critical point is not degenerate. By a straightforward computation, starting by differentiating the equation  $\partial_{\vartheta}^2 \kappa = 1/\kappa - \kappa$ , one gets

$$\begin{split} \partial_{\vartheta}^{3}\kappa^{2} &= 2\partial_{\vartheta}^{2}(\kappa\partial_{\vartheta}\kappa) = 2\partial_{\vartheta}(\partial_{\vartheta}\kappa)^{2} + 2\partial_{\vartheta}(\kappa\partial_{\vartheta}^{2}\kappa) = 6\partial_{\vartheta}\kappa\partial_{\vartheta}^{2}\kappa + 2\kappa\partial_{\vartheta}\partial_{\vartheta}^{2}\kappa \\ &= 6\frac{\partial_{\vartheta}\kappa}{\kappa} - 6\kappa\partial_{\vartheta}\kappa - 2\frac{\kappa}{\kappa^{2}}\partial_{\vartheta}\kappa - 2\kappa\partial_{\vartheta}\kappa = 4\frac{\partial_{\vartheta}\kappa}{\kappa} - 4\partial_{\vartheta}\kappa^{2} \end{split}$$

so that

$$\partial_{\vartheta}^{3} \kappa^{2} + 4 \partial_{\vartheta} \kappa^{2} = 4 \frac{\partial_{\vartheta} \kappa}{\kappa} .$$

We compute

$$4 \int_{0}^{\vartheta_{0}} \sin(2\vartheta) \frac{\partial_{\vartheta} \kappa}{\kappa} d\vartheta = \int_{0}^{\vartheta_{0}} \sin(2\vartheta) \left( \partial_{\vartheta}^{3} \kappa^{2} + 4 \partial_{\vartheta} \kappa^{2} \right) d\vartheta$$

$$= \sin(2\vartheta) \partial_{\vartheta}^{2} \kappa^{2} \Big|_{0}^{\vartheta_{0}} - 2 \int_{0}^{\vartheta_{0}} \cos(2\vartheta) \partial_{\vartheta}^{2} \kappa^{2} d\vartheta + 4 \int_{0}^{\vartheta_{0}} \sin(2\vartheta) \partial_{\vartheta} \kappa^{2} d\vartheta$$

$$= 2 \sin(2\vartheta_{0}) \left( \kappa(\vartheta_{0}) \partial_{\vartheta}^{2} \kappa(\vartheta_{0}) + (\partial_{\vartheta} \kappa)^{2} (\vartheta_{0}) \right) - 2 \cos(2\vartheta) \partial_{\vartheta} \kappa^{2} \Big|_{0}^{\vartheta_{0}}$$

$$- 4 \int_{0}^{\vartheta_{0}} \sin(2\vartheta) \partial_{\vartheta} \kappa^{2} d\vartheta + 4 \int_{0}^{\vartheta_{0}} \sin(2\vartheta) \partial_{\vartheta} \kappa^{2} d\vartheta$$

$$= 2 \sin(2\vartheta_{0}) \left( \kappa(\vartheta_{0}) \partial_{\vartheta}^{2} \kappa(\vartheta_{0}) + (\partial_{\vartheta} \kappa)^{2} (\vartheta_{0}) \right)$$

$$- 4 \cos(2\vartheta_{0}) \kappa(\vartheta_{0}) \partial_{\vartheta} \kappa(\vartheta_{0}) + 4 \kappa(0) \partial_{\vartheta} \kappa(0) .$$

Now, since  $\partial_{\vartheta}\kappa(0) = \partial_{\vartheta}\kappa(\vartheta_0) = 0$  using the equation for the curvature  $\partial_{\vartheta}^2\kappa = 1/\kappa - \kappa$  we get

$$4\int_0^{\vartheta_0} \sin(2\vartheta) \frac{\partial_{\vartheta} \kappa}{\kappa} d\vartheta = 2\sin(2\vartheta_0)(1 - \kappa^2(\vartheta_0)),$$

and this last term is nonnegative as  $\kappa < 1$  at a local minimum and  $0 < 2\vartheta_0 \le \pi$ . Looking at the left-hand integral we see instead that the factor  $\sin(2\vartheta)$  is always nonnegative, since  $2\vartheta_0 \le \pi$  and  $\partial_{\vartheta}\kappa$  is always nonpositive in the interval  $[0, \vartheta_0]$ ,

as we assumed that we were moving from a local maximum of  $\kappa$  at 0 to a local minimum of  $\kappa$  at  $\vartheta_0$  without crossing any other critical point of  $\kappa$ . This gives a contradiction so  $\Sigma$  must be the unit circle.

## 3. Convex hypersurfaces with pinched second fundamental form

**Definition 3.1** (Complete Riemannian manifold). A (geodesically) complete manifold is a Riemannian manifold for which every maximal (inextendible) geodesic is defined on  $\mathbb{R}$ .

**Definition 3.2** (Conformal map). Two maps  $X,Y:M^n\to\mathbb{R}^{n+1}$  are *conformal*, if there exists  $\lambda:M^n\to\mathbb{R}$  with

$$g_{ij}^X = \lambda g_{ij}^Y$$
.

We say X is quasi-conformal with respect to Y if

$$g_{ij}^X \ge \lambda \, g_{ij}^Y$$
.

See [Ham94]. Suppose that  $M = X(M^n) \subset \mathbb{R}^{n+1} = \mathbb{R}^n \times \mathbb{R}$  is written as a graph over a convex over a convex open set  $U \subset \mathbb{R}^n$  of a strictly convex function

$$y = f(x_1, \ldots, x_n)$$

so that  $y \to \infty$  as  $x = (x_1, \dots, x_n) \to \partial U$ . By translating upwards if necessary, since y is bounded below, we can assume  $y \ge e$  everywhere, so that  $\log \log y \ge 0$ . Let  $g_{ij}$  be the Riemannian metric induced on M so that

$$g_{ij} = \delta_{ij} + \frac{\partial y}{\partial x^i} \frac{\partial y}{\partial x^j} .$$

**Theorem 3.3** (Hamilton, [Ham94, Theorem 2.1]). The conformally equivalent metric

$$\tilde{g}_{ij} = \frac{g_{ij}}{(y \log y)^2}$$

is complete with finite volume.

*Proof.* First, we show that  $\tilde{g}_{ij}$  is complete. We have  $\det(g_{ij}) \geq 1$ . For any geodesic  $\gamma: I \to M$  going to infinity, we have  $\gamma^n \to \infty$ . Therefore its length satisfies,

$$\begin{split} \tilde{L}(\gamma) &= \int_{I} |\gamma'(t)|_{\tilde{g}} \, dt \geq \int_{\gamma^{n}(a)}^{\infty} \sqrt{\det(\tilde{g}_{ij})} dy \\ &\geq \int_{\gamma^{n}(a)}^{\infty} \frac{dy}{y \log y} = \log \log y|_{\gamma^{n}(a)}^{\infty} = \infty \,. \end{split}$$

Since geodesics have constant speed, this is what we desired. To estimate the volume, we observe that, because y is a strictly convex function of x, outside a compact set we must have

$$\left| \frac{\partial y}{\partial x^i} \right| \ge \delta$$

for some  $\delta > 0$  and at least one  $i \in \{1, ..., n\}$ . Let dV denote the volume element on M in the induced metric  $g_{ij}$ , which in x coordinates is

$$dV = \sqrt{\det\left(\delta_{ij} + \frac{\partial y}{\partial x^i} \frac{\partial y}{\partial x^j}\right)} dx^1 \dots dx^n.$$

Let  $k \in \mathbb{N}$  and

$$M^k := M \cap \{e + k - 1 \le y \le e + k\}$$

and let  $dV^k$  denote the volume element of the part of  $M^k$ . We can devide  $M^k$  into pieces  $M_1^k, \ldots, M_n^k$ , where  $\frac{\partial y}{\partial x^i}$  is largest on  $M_i^k$ , and estimate  $dV_i^k$  from above

on each piece. For each  $k \in \mathbb{N}$ , on  $M_i^k$ , we take  $x^1, \dots, x^{i-1}, y, x^{i+1}, \dots, x^n$  as coordinates. Since  $\frac{\partial y}{\partial x^i}$  is larger than the other derivatives, and  $\left|\frac{\partial y}{\partial x^i}\right| \geq \delta > 0$ ,

$$\sqrt{\det\left(\delta_{ij} + \frac{\partial y}{\partial x^i} \frac{\partial y}{\partial x^j}\right)} \le C \left| \frac{\partial y}{\partial x^i} \right|$$

and thus

$$dV_i^k \le C dx^1 \dots dx^{i-1} dy dx^{i+1} \dots dx^n$$

on  $M_i^k$ . By the gradient estimate shows that

$$|x| \le Cy$$

for a suitable large constant. Let

$$U_i^k := \left\{ x \in \mathbb{R}^n \mid (x, f(x)) \in M_i^k \right\}.$$

We can integrate in every direction  $x^1, \ldots, x^{i-1}, x^{i+1}, \ldots x^n$  and estmate

$$\int_{U_i^k} dV_i^k \le C \int_{U_i^k} dx^1 \dots dx^{i-1} dy dx^{i+1} \dots dx^n \le C \int_{U_i^k} y^{n-1} dy,$$

that is,

$$dV_y^i \le Cy^{n-1}dy$$

Hence,

$$d\tilde{V}_y^i \le \frac{Cdy}{y \log^n y}$$

and

$$\tilde{V} = \int_{U} d\tilde{V} = \sum_{k \in \mathbb{N}} \sum_{i=1}^{n} \int_{U_{i}^{k}} d\tilde{V}_{i}^{k} \le C \sum_{k \in \mathbb{N}} \sum_{i=1}^{n} \int_{U_{i}^{k}} \frac{dy}{y \log^{n} y}$$

$$= C \int_{e}^{\infty} \frac{dy}{y \log^{n} y} = \frac{-C}{(n-1) \log^{n-1} y} \Big|_{e}^{\infty} = \frac{C}{n-1} < \infty.$$

- **Remark 3.4.** (i) Let X be an embedding of the  $\mathbb{S}^{n-1}$ , Y the standard embedding of the  $\mathbb{S}^n$ . Let  $p_N, p_S \in \mathbb{S}^n$  and  $N, S \in \mathbb{S}^n$  the north and south pole. Show: there exists a conformal map Define  $Z: \mathbb{S}^n \to \mathbb{S}^n$  with  $Z(p_N) = p_N$  and  $Z(p_S) = S$ .
- (ii) Let X be an embedding of the  $\mathbb{S}^{n-1}$ , Y be an embedding of the  $\mathbb{S}^n$  and Z be an embedding of the cylinder  $\mathbb{S}^{n-1} \times [-R, R]$ , where

$$Y(x, \vartheta) = (X(x)\cos(\vartheta), \sin(\vartheta))$$

and

$$Z(x, \vartheta) = (X(x), z(\vartheta))$$

for  $\vartheta \in [-\pi/2, \pi/2)$ . Then

$$(g_{ij}^Y) = \begin{pmatrix} \cos^2(\vartheta) \, g_{ij}^X & 0\\ 0 & 1 \end{pmatrix}$$

and

$$(g_{ij}^Z) = \begin{pmatrix} g_{ij}^X & 0 \\ 0 & (z'(\vartheta))^2 \end{pmatrix} .$$

For Y and Z to be conformal with  $(g_{ij}^Y) = \lambda(g_{ij}^Z)$ , we have to choose

$$\lambda(\vartheta) = \cos^2(\vartheta)$$
 and  $z'(\vartheta) = \frac{1}{\cos(\vartheta)}$ 

for  $\vartheta \in [-\pi/2 + \varepsilon, \pi/2 - \varepsilon]$ , where  $\varepsilon > 0$  and  $R = R(\varepsilon)$ , which is realized by

$$z(\vartheta) = \log\left(\tan\left(\frac{\theta}{2} + \frac{\pi}{4}\right)\right) \,.$$

**Theorem 3.5** (Hamilton, [Ham94]). Let U be an open subset of the unit sphere  $\mathbb{S}^n$  which is not empty and whose closure is not the whole sphere. Then there is no metric on U, conformal with respect to the round metric, which is complete with finite volume.

*Proof.* By hypotheses we can find some point  $p_N$  which is contained in U, and some point  $p_S$  which avoids the closure of U. By Remark 3.4, we can assume that  $p_N$  is the north pole and  $p_S$  is the south pole. We can then find an  $\varepsilon > 0$  so that the  $\varepsilon$ -ball around  $p_N$  lies in U,

$$B_{\varepsilon}(p_N) \subset U$$

while the  $\varepsilon$ -ball around  $p_S$  avoids U,

$$B_{\varepsilon}(p_S) \subset \mathbb{S}^n \setminus U$$
.

By Remark 3.4, we can find a conformal map  $\varphi$  of the sphere  $\mathbb{S}^n$  minus these two balls to the cylinder  $\mathbb{S}^{n-1} \times [0, L]$ ,

$$\varphi: \mathbb{S}^n \to \mathbb{S}^{n-1} \times [0, L]$$

taking the boundary of the  $\varepsilon$ -ball around  $p_N$  to  $\mathbb{S}^{n-1} \times \{0\}$ 

$$\varphi(\partial B_{\varepsilon}(p_N)) = \mathbb{S}^{n-1} \times \{0\}$$

and the boundary of the  $\varepsilon$ -ball around  $p_S$  to  $\mathbb{S}^{n-1} \times \{L\}$ ,

$$\varphi(\partial B_{\varepsilon}(p_S)) = \mathbb{S}^{n-1} \times \{L\}.$$

The part of U outside the  $\varepsilon$ -ball around  $p_N$  will map to some relatively open subset

$$W := \varphi(U \setminus B_{\varepsilon}(p_N)) \subset (\mathbb{S}^{n-1} \times [0, L]) \setminus (\mathbb{S}^{n-1} \times \{L\})$$

of the cylinder which contains  $\mathbb{S}^{n-1} \times \{0\}$  and avoids  $\mathbb{S}^{n-1} \times \{L\}$ ,

$$\mathbb{S}^{n-1}\times\{0\}\subset W\,.$$

The subset W will be a noncompact manifold with one compact boundary component  $\mathbb{S}^{n-1}$ . Any complete metric

$$q^U$$
 on  $U$ 

with finite volume conformal to the round metric

$$a^{\mathbb{S}^n}$$
 on  $\mathbb{S}^n$ 

would give a complete metric with finite volume on

$$g^W$$
 on  $W$ 

conformal to the product metric

$$g^{\mathbb{S}^{n-1}\times[0,L]}$$
 on  $\mathbb{S}^{n-1}\times[0,L]$ .

We show that such cannot exist. We introduce coordinates

$$\vartheta = (\vartheta_1, \dots, \vartheta_{n-1})$$
 on  $\mathbb{S}^{n-1}$ 

and

t on 
$$[0,L]$$

Let  $g^{\mathbb{S}^{n-1}}$  denote the metric on  $\mathbb{S}^{n-1}$  and  $d\mu$  the volume form. Then

$$g := g^{\mathbb{S}^{n-1} \times [0,L]} = \begin{pmatrix} g^{\mathbb{S}^{n-1}} & 0 \\ 0 & 1 \end{pmatrix}$$

is the product metric on  $\mathbb{S}^{n-1} \times [0, L]$  and

$$dV = d\mu dt$$

is the product volume form. For every  $\vartheta \in \mathbb{S}^{n-1}$ , there will be a first point

$$t = h(\vartheta)$$

where the pair  $(\vartheta,t)$  is no longer in W. Of course h may not be a continuous function and the pair may reenter W for larger values of t. This does not matter. Any quasi-conformally equivalent metric on W is given by

$$\tilde{g} = \lambda(\vartheta, t)g$$

for some funtion  $\lambda$  defined at least for  $0 \leq t \leq h(\vartheta)$ . The corresponding volume form is

$$d\tilde{V} = \lambda^n d\mu dt.$$

If the total volume  $\tilde{V}$  of W in the conformally equivalent metric is finite, we have

$$\iint_{W} \lambda^{n} d\mu dt = \tilde{V} < \infty.$$

By Hölder's inequality

$$\iint_{W} \lambda \, d\mu dt \le \left( \iint_{W} \lambda^{n} d\mu dt \right)^{1/n} \left( \iint_{W} d\mu dt \right)^{(n-1)/n}$$

and surely

$$\iint_W d\mu dt \le L|\mathbb{S}^{n-1}| < \infty.$$

Therefore

$$\iint_{0\leq t< h(\vartheta)} \lambda(\vartheta,t) d\mu dt <\infty\,.$$
 On the other hand, if we integrate first in  $t$ , we see that

$$\int_{\mathbb{S}^{n-1}} \left( \int_0^{h(\vartheta)} \lambda(\vartheta,t) \, dt \right) d\mu \geq |\mathbb{S}^{n-1}| \inf_{\vartheta \in \mathbb{S}^{n-1}} \int_0^{h(\vartheta)} \lambda(\vartheta,t) \, dt$$

and therefore

$$\inf_{\vartheta \in \mathbb{S}^{n-1}} \int_0^{h(\vartheta)} \lambda(\vartheta,t) \, dt < \infty \, .$$

But along a path where  $\vartheta$  is constant we have  $\tilde{q} = \lambda$ . Thus there is some  $\vartheta$  where the path from  $(\vartheta, 0)$  to  $(\vartheta, h(\vartheta))$  has finite length. This shows that the metric is not complete and proves the theorem.

**Theorem 3.6** (Hamilton, [Ham94, Theorem 1.1]). Let M be a smooth strictly convex hypersurface bounding a region in  $\mathbb{R}^{n+1}$ ,  $n \geq 2$ . Suppose that its second fundamental form is  $\varepsilon$ -pinched in the sense that

$$h_{ij} \geq \varepsilon H g_{ij}$$

for some  $\varepsilon > 0$ . Then M is compact.

*Proof.* Assume that M is noncompact. By Theorem 3.3, M has a conformally equivalent metric  $\tilde{g}_{ij}$  which is complete with finite volume. Observe that the Gauss map  $\nu: M \to \mathbb{S}^n$  gives a diffeomorphism of the convex hypersurface M onto an open subset  $U = \nu(M)$  of the sphere  $\mathbb{S}^n$  which lies in a hemisphere. Thus U is not empty and its closure is not all of  $\mathbb{S}^n$ . By Theorem 3.5, there is no metric  $\hat{g}_{ij}$  on U, quasi-conformal with respect to the round metric, which is complete with finite volume. However, the pinching condition implies

$$\varepsilon H \delta_i^k \leq h_i^k \leq H \delta_i^k$$

so that

$$\varepsilon H \partial_i = \varepsilon H \delta_i^k \partial_k \le h_i^k \partial_k = \partial_i \nu \le H \delta_i^k \partial_k = H \partial_i$$
.

We define

$$\hat{g}_{ij} := \langle \partial_i \boldsymbol{\nu}, \partial_j \boldsymbol{\nu} \rangle$$

and observe that

$$(\varepsilon H)^2 g_{ij} = (\varepsilon \tilde{H})^2 \tilde{g}_{ij}.$$

Hence,

$$(\varepsilon \tilde{H})^2 \tilde{g}_{ij} \le \hat{g}_{ij} \le \tilde{H}^2 \tilde{g}_{ij}$$
.

If  $\tilde{g}_{ij}$  is complete,  $\hat{g}_{ij}$  is, by the first inequality. If  $\tilde{g}_{ij}$  has finite Volume,  $\hat{g}_{ij}$  must have by the second inequality. This is a contradiction.

#### 4. Singularities

**Definition 4.1** (Singularities, see [Eck04, Definitions 3.5 and 5.1]). We say that a solution  $(M_t)_{t\in[0,T)}$  of (MCF) reaches a point  $x_0 \in \mathbb{R}^{n+1}$  at time  $T \leq \infty$  if there exists a sequence  $(p_k, t_k)_{k\in\mathbb{N}}$  in  $M^n \times [0, T)$  with  $t_k \nearrow T$  so that  $X(p_k, t_k) \to x_0$  for  $k \to \infty$ .

Let  $\mathcal{S}$  be the set of points  $x \in \mathbb{R}^{n+1}$  so that there exists a sequence  $(p_k, t_k)_{k \in \mathbb{N}}$  with  $t_k \nearrow T$  and  $X(p_k, t_k) \to x$  for  $k \to \infty$ . We call  $\mathcal{S}$  the set of reachable points.

A point  $x_0 \in \mathbb{R}^2$  is called a *singular* or *blow-up point* of the flow at time T if  $(M_t)_{t\in[0,T)}$  reaches  $x_0$  at time T and has no smooth extension beyond time T in any neighbourhood of  $x_0$ . The sequence  $(p_k, t_k)_{k\in\mathbb{N}}$  is called *blow-up sequence*.

All other points (which includes those not reached by the solution) are called regular points.

We want to investigate singularities of the flow.

**Proposition 4.2.** Let  $T < \infty$ . If  $|A|^2 \le C_0$  on  $M^n \times [0,T)$ , then  $|\nabla^m A|^2 \le C_m$  on  $M^n \times [0,T)$ , where  $C_m = C_m(n, M_0, C_0)$ .

Proof. See [Sch17c, Proposition 2.1.5]. By Lemma 1.4,

$$\partial_t |\nabla^m A|^2 \leq \Delta |\nabla^m A|^2 - 2|\nabla^{m+1} A|^2 + C(n,m) \sum_{i+j+k=m} |\nabla^i A| |\nabla^j A| |\nabla^k A| |\nabla^m A| \,.$$

We give a proof by induction. The case m=0 is trivially true. So we assume that for m>0 we have  $|\nabla^l A|^2 \leq C_l$  for  $0\leq l\leq m-1$ . Thus

$$\partial_t |\nabla^{m-1} A|^2 \le \Delta |\nabla^{m-1} A|^2 - 2|\nabla^m A|^2 + B_{m-1}$$

and

$$\partial_t |\nabla^m A|^2 \le \Delta |\nabla^m A|^2 - B_m \left(1 + |\nabla^m A|^2\right).$$

We consider the function  $f := |\nabla^m A|^2 + B_m |\nabla^{m-1} A|^2$ , which satisfies

$$\partial_t f \le \Delta f - B_m \left( 1 + |\nabla^m A|^2 \right) - 2B_m |\nabla^m A|^2 + B_{m-1} B_m$$

$$\le \Delta f - B_m f + B_m^2 |\nabla^{m-1} A|^2 + B_{m-1} B_m$$

$$\le \Delta f - B_m f + B.$$

Define  $\tilde{f} := \exp(B_m t) f - \exp(B_m T) B t$ . Then

$$\partial_t \tilde{f} \le \exp(B_m t)(B_m f + \partial_t f) - \exp(B_m T) B$$
  
 $\le \exp(B_m t)(\Delta f + B) - \exp(B_m T) B \le \Delta \tilde{f}$ 

which implies  $\tilde{f}(\cdot,t) \leq \max_{M} \tilde{f}(\cdot,0)$  and thus

$$f(\cdot,t) \le \exp(-B_m t) \left(\max_{M} \tilde{f}(\cdot,0) + \exp(B_m T) B t\right) \le C.$$

**Theorem 4.3.** Let  $T < \infty$  and  $(M_t)_{t \in [0,T)}$  be a family of smooth, immersed hypersurfaces evolving by (MCF) with

$$M_t \cap B_R(0) \neq \emptyset$$

for some R > 0 and all  $t \in [0,T)$  and there exists  $C_0 < \infty$  such that

$$\sup_{t \in [0,T)} \sup_{M_t} |A| \le C_0.$$

Then  $M_T$  is smooth.

*Proof.* By Proposition 4.2,

$$\sup_{t \in [0,T)} \sup_{M_t} |\nabla^m|A|| \le C_m$$

for all  $m \in \mathbb{N} \cup \{0\}$ . By Lemma 1.4,

$$\partial_t \boldsymbol{\nu} = \nabla H$$

so that the rotation of the normal is uniformly bounded in small space-time neighbourhoods. That is, there exist  $t_0 \in [0,T)$ , r > 0 and  $\varepsilon > 0$  so that for each  $p \in M^n$  there exists an open neighbourhood

$$U_{r,t_0}(p) = X^{-1}(B_r(X(p,t_0)),t_0) \subset \mathbb{R}^n$$
,

where  $B_r$  is the geodesic ball in  $M_{t_0}$ , so that, after rotation and translation,

$$\nu(q,t) \in \mathbb{S}^n \cap \{x^n \ge \varepsilon\}$$

for all  $q \in U_{r,t_0}(p)$  and  $t \in [t_0,T)$ . For  $R_0 \geq R$ , there exist finitely many points  $\{p_i\}_{i=1}^{N_0}$  so that

$$M_t \cap B_{R_0}(0) \subset \bigcup_{i=1}^{N_0} X(U_{r,t_0}(p_i),t)$$

for all  $t \in [t_0, T)$ . For  $p \in \{p_i\}_{i=1}^{N_0}$  we can write

$$M_t \cap X(U_{r,t_0}(p),t)$$

as a graph of a function  $f: U_{r,t_0}(p) \times [t_0,T) \to \mathbb{R}$  with  $|D^m f|$  uniformly bounded on  $U_{r,t_0}(p) \times [t_0,T)$  for all  $m \in \mathbb{N} \cap \{0\}$ . Let  $(t_k)_{k \in \mathbb{N}}$  with  $t_k \nearrow T$ . By Arzelá–Ascoli, for each  $m \in \mathbb{N} \cap \{0\}$ , the sequence

$$(f_k^m := D^m f(\cdot, t_k))_{k \in \mathbb{N}}$$

converges uniformly along a subsequence to a continuous limit

$$f_{\infty}^m = D^m f_{\infty} = D^m f(\cdot, T)$$
.

Hence,  $f(\cdot,T)$  is smooth. This can be done for each  $i \in \{1,\ldots,N_0\}$ . We define

$$X_k := X(\cdot, t_k)$$
.

Locally, we can describe  $X_k$  via  $f_k$ . Thus  $X(\cdot,T)$  is smooth on  $\bigcup_{i=1}^{N_0} U_{r,t_0}(p_i)$  and so is  $M_T \cap B_{R_0}(0)$ . Let now be  $(R_l)_{l \in \mathbb{N}}$  be a sequence of radii with  $R \leq R_l \nearrow \infty$ . For each  $l \in \mathbb{N}$ , there exist finitely many points  $\{p_i\}_{i=1}^{N_l}$  so that

$$M_t \cap B_{R_l}(0) \subset \bigcup_{i=1}^{N_l} X(U_{r,t_0}(p_i),t)$$

for all  $t \in [t_0, T)$ . Define

$$X_k^l := X^l(\cdot, t_k)$$

locally via  $f_k^l$ . By the same argument as above,  $X_{\infty}^l = X^l(\cdot,T): \bigcup_{i=1}^{N_l} U_{r,t_0}(p_i) \to \mathbb{R}^{n+1}$  and  $M_T \cap B_{R_l}(0)$  is smooth for every  $l \in \mathbb{N}$ . We now pick a diagonal sequence to obtain a smooth limit  $X_{\infty}^{\infty} = X(\cdot,T): M^n \to \mathbb{R}^{n+1}$  with image  $M_T$  which coincides with  $X_{\infty}^l$  on every ball  $B_{R_l}(0)$ . Since  $M_t \to M_T$  continuously for  $t \to T$ , the smooth convergence holds for  $t \to T$ .

Corollary 4.4. If  $T < \infty$ , then  $\limsup_{t \to T} \max_{M_t} |A|^2 = \infty$ .

Proof. See [Sch17c, Corollary 2.1.6]. Let us assume to the contrary that  $|A|^2 \leq C_0$  on  $M^n \times [0,T)$ . By Proposition 4.2 all higher derivatives of A are uniformly bounded on  $M^n \times [0,T)$ . By Theorem 4.3,  $X(\cdot,T)$  is a smooth immersion. By short-time existence this implies that we can extend the solution further, which contradicts the assumption that T is maximal.

**Lemma 4.5** (Hamilton's trick [Ham86, Lemma 3.5]). Let  $f : [a,b] \times (0,T) \to \mathbb{R}$  be in  $C^1$ . Then  $f_{\max}(t) := \max_{p \in [a,b]} f(p,t)$  is locally Lipschitz for  $t \in (0,T)$  and at a differentiable time,

$$\frac{d}{dt}f_{\max}(t) \le \sup \left\{ \partial_t f(p,t) \mid p \in [a,b] \text{ with } f(p,t) = f_{\max}(t) \right\}.$$

**Proposition 4.6** (Huisken, [Hui90, Lemma 1.2]). If  $T < \infty$ , then max  $|A|^2(t) \to \infty$  for  $t \to T$  where

$$\max |A|^2(t) \ge \frac{1}{\sqrt{2(T-t)}}.$$

*Proof.* By Corollary 4.4,  $|A|_{\max}(t) \to \infty$  for  $t \to T$ . For  $t \in (0,T)$ , let  $p \in M^n$  so that  $|A|^2(p,t) = |A|^2_{\max}(t)$ . Then

$$\operatorname{Hess} |A|^2(p,t) \leq 0$$
.

By Lemma 1.4

$$\partial_t |A|^2 = \Delta |A|^2 - |\nabla A|^2 + 2|A|^4 \le 2|A|^4$$

at (p,t). Since  $|A|_{\max}^2$  is Lipschitz, by Rademacher's theorem, Theorem A.3,  $\partial_t |A|_{\max}^2$  exists for almost every  $t \in (0,T)$ . By Hamilton's trick, Lemma 4.5,

$$\begin{aligned} \partial_t |A|_{\max}^2(t) &\leq \max \left\{ \partial_t |A|^2(p,t) \mid p \in M^n \text{ with } |A|^2(p,t) = |A|_{\max}^2(t) \right\} \\ &\leq \max \left\{ 2|A|^4(p,t) \mid p \in M^n \text{ with } |A|^2(p,t) = |A|_{\max}^2(t) \right\} = 2|A|_{\max}^4(t) \end{aligned}$$

for almost every  $t \in (0,T)$ . Assume that there exists a time  $t_0 \in [0,T)$  where  $|A|_{\max}^2 = 0$ . Then  $M_{t_0}$  is a plane segment in  $\mathbb{R}^{n+1}$  which contradicts that  $T < \infty$ . Hence,  $|A|_{\max}^2(t) > 0$  for all  $t \in [0,T)$  and  $|A|_{\max}^{-2}$  is Lipschitz as well. Rademacher's theorem implies that  $\partial_t |A|_{\max}^{-2}(t)$  exists for almost every  $t \in (0,T)$ . Thus,

$$\partial_t |A|_{\max}^{-2} = -|A|_{\max}^{-4} \partial_t |A|_{\max}^2 \ge -2$$
 (7)

for almost every  $t \in (0,T)$ . Since  $|A|_{\max}^{-2}$  is Lipschitz, we can integrate (7) over an interval  $[t,t_k] \subset [0,T)$  to obtain

$$\frac{1}{|A|_{\max}^2(t_k)} - \frac{1}{|A|_{\max}^2(t)} \ge -2(t_k - t). \tag{8}$$

Let  $t \in [0,T)$  and  $(t_k)_{k \in \mathbb{N}}$  be a sequence with  $t_k \in (t,T)$  for all  $k \in \mathbb{N}$ ,  $t_k \nearrow T$  and  $|A|^2_{\max}(t_k) \to \infty$  for  $k \to \infty$ . Taking the limit  $k \to \infty$  in (8) yields

$$\frac{1}{|A|_{\max}^2(t)} \le 2(T-t)$$

for all  $t \in [0,T)$ .

**Example 4.7.** (i) The curvature of the spheres  $\mathbb{S}_{r(t)}^n$  blows up in the exact rate. (ii) A dumbbell with a small neck develops a singularity at the neck before the surface disappears.

We distinguish between two types of singularities.

**Definition 4.8** (Type-I and type-II singularities). We say that a singularity is of type I, if there exists a constant  $C_0 > 1$  so that

$$|A|_{\max}(t) \le \frac{C_0}{\sqrt{T-t}} \tag{9}$$

for all  $t \in [0, T)$ , and of type II, if such a constant does not exist, that is,

$$\limsup_{t \to T} |A|_{\max}(t)\sqrt{T - t} = \infty.$$
 (10)

**Remark 4.9** (Parabolic rescaling). Let  $\lambda > 0$  and  $t_0 \in (0,T)$ . Consider the rescaled flow  $X_{\lambda}: M^n \times [-\lambda^2 t_0, t_0) \to \mathbb{R}^2$  with

$$X_{\lambda}(p,\tau) = \lambda \left( X \left( p, t_0 + \frac{\tau}{\lambda^2} \right) - x_0 \right).$$

and define

$$M_{\tau}^{\lambda} := \lambda \left( M_{t_0 - \tau/\lambda^2} - x_0 \right) .$$

Then  $\tau = \lambda^2(t-t_0)$ ,  $\partial_{\tau} = \frac{1}{\lambda^2}\partial_t$ ,  $g_{ij}^{\lambda} = \lambda^2 g_{ij}$  and  $h_{ij}^{\lambda} = \lambda h_{ij}$  so that

$$|A_{\lambda}| = \frac{1}{\lambda}|A|$$
 and  $H_{\lambda} = \frac{1}{\lambda}H$ 

so that

$$\partial_{\tau}X_{\lambda} = \frac{1}{\lambda}\partial_{t}X = -\frac{1}{\lambda}H\boldsymbol{\nu} = -H_{\lambda}\boldsymbol{\nu}$$

again flows by mean curvature flow.

**Theorem 4.10.** Let  $T < \infty$  and  $k \in \mathbb{N}$ . Let  $\emptyset \neq J_k \subset J_{k+1}$  be a sequence of intervals and  $(M_{\tau}^k)_{\tau \in J_k}$  be families of smooth, immersed hypersurfaces evolving by (MCF) for each  $k \in \mathbb{N}$  with

$$M_{\tau}^k \cap B_R(0) \neq \emptyset$$

for some R > 0 and for all  $k \in \mathbb{N}$  and all  $\tau \in J_k$ , and there exists  $C_0 < \infty$  such that

$$\sup_{k \in \mathbb{N}} \sup_{\tau \in J_k} \sup_{M_{\tau}^k} |A_k| \le C_0.$$

Then there exists a subsequence  $((M_{\tau}^k)_{\tau \in J_k})_{k \in \mathbb{N}}$  that converges on compact subsets of  $J_{\infty}$  and in  $\mathbb{R}^{n+1}$  to a smooth, immersed limit flow  $(M_{\tau}^{\infty})_{\tau \in J_{\infty}}$  evolving by (MCF).

*Proof.* By Proposition 4.2,

$$\sup_{k \in \mathbb{N}} \sup_{\tau \in J_k} \sup_{M_{\tau}^k} |\nabla^m| A_k || \le C_m$$

for all  $m \in \mathbb{N} \cup \{0\}$ . Let  $R_0 \geq R$ ,  $k_0 \in \mathbb{N}$  and  $\tau_0 \in J_k$  for  $k \geq k_0$ . Since  $M_{\tau_0}^k$  is smooth and

$$\tilde{M}_{\tau_0}^k := M_{\tau_0}^k \cap B_{R_0}^{n+1}(0) \neq \emptyset$$

for every  $k \in \mathbb{N}$ , there exists a subsequence  $(\tilde{M}_{\tau_0}^k)_{k \in \mathbb{N}}$  with continuous limit

$$\tilde{M}_{\tau_0}^{\infty} \subset B_{R_0}^{n+1}(0)$$
.

Moreover, there exists r > 0 so that for every  $x \in \tilde{M}_{\tau_0}^{\infty}$ ,

$$\tilde{M}_{\tau_0,r}^{\infty}(x) := \tilde{M}_{t_0}^{\infty} \cap B_r^{n+1}(x)$$

can be written as a graph of some function  $g: B^n_r(x) \subset P(x) \to \mathbb{R}$  over some affine tangent plane P(x) at x. By the convergence, there exists a subsequence  $(\tilde{M}^k_{t_0})_{k \in \mathbb{N}}$  so that, for k big enough,

$$\tilde{M}^k_{\tau_0}\cap B^{n+1}_r(x)$$

can be written as graphs of some function  $g_k: B^n_{r/2}(x) \to \mathbb{R}$  over the same affine plane P(x). By the uniform bounds on  $|A_k|$ ,  $|D^m g_k|$  is uniformly bounded for all  $m \in \mathbb{N}$  and  $g_k$  is smooth for every  $k \geq k_0$ . Furthermore, there exists  $\delta, \varepsilon > 0$  so that, after rotation and translation,

$$\nu_k(y) \in \mathbb{S}^n \cap \{x^n > \varepsilon\}$$

for all  $y \in \tilde{M}_{\tau}^k \cap B_r^{n+1}(x)$  and  $\tau \in (\tau_0 - \delta, \tau_0 + \delta)$ , so that  $\tilde{M}_{\tau}^k \cap B_r^{n+1}(x)$  can be written as graphs of the functions  $f_k : B_{r/2}^n(x) \times (t_0 - \delta, t_0 + \delta) \to \mathbb{R}$ . Since all time derivatives can be expressed in terms of spatial derivatives, f is smooth in time. By Arzelá–Acsoli,  $(f_k)_{k \in \mathbb{N}}$  converges along a subsequence to a smooth limit  $f_{\infty}$ . Like in the proof of Theorem 4.3, we can repeat this process this for a sequence  $(R_l)_{l \in \mathbb{N}}$ 

with  $R \geq R_l \to \infty$ , and after picking a diagonal sequence we obtain a smooth limit  $M_{\tau}^{\infty} \subset \mathbb{R}^{n+1}$ . Note that a subsequence of the  $X_k(\cdot, \tau)$  does not necessarily converge to a limiting immersion; it will be necessary to "reparametrize"  $X_k(\cdot, \tau)$  (see [Lan85, ] for details).

#### 5. Typ-I singularities

We want to rescale the surface  $M_t$  near a type-I singularity as  $t \to T < \infty$ . The following rescaling technique was introduced in [HS99, Remark 4.6].

**Definition 5.1** (Type-I rescaling). Let  $(p_k, t_k)_{k \in \mathbb{N}}$  be a blow-up sequence in  $M^n \times [0, T)$  with  $t_k \nearrow T$  for  $k \to \infty$  and

$$|A|^2(p_k, t_k) = \max_{p \in M^n} |A|^2(p, t_k) = \max_{M^n \times [0, t_k]} |A|^2(p, t)$$

for each  $k \in \mathbb{N}$ . We set

$$\lambda_k^2 := |A|^2(p_k, t_k)$$
 and  $\alpha_k := -\lambda_k^2 T$ 

and define the rescaled embeddings  $X_k:M^n\times [\alpha_k,0)\to \mathbb{R}^2$  by

$$X_k(p,\tau) := \lambda_k \left( X \left( p, T + \frac{\tau}{\lambda_k^2} \right) - x_0 \right). \tag{11}$$

**Lemma 5.2** (Properties of the type-I rescaling). Let  $X: M^n \times (0,T) \to \mathbb{R}^2$  be a solution of (MCF) with  $T < \infty$ . For the type-I rescaling 5.1 in case of a type-I singularity,

$$\lambda_k \to \infty$$
 and  $\alpha_k \to -\infty$ 

for  $k \to \infty$ . Furthermore,

$$X_k(0,\tau_k) \in B_{3C_0^2}(0)$$
 and  $|A_k|^2(0,\tau_k) = 1$ ,

where

$$\tau_k := -\lambda_k^2 (T - t_k) \in \left[ -\frac{C_0^2}{2}, -\frac{1}{2} \right]$$

and, for  $\delta > 0$ ,

$$\max_{M^n \times [\alpha_k, -\delta^2]} |A_k| \le \frac{C_0}{\delta}$$

for all  $k \in \mathbb{N}$ .

*Proof.* We follow [MB14, Corollary 4.8, Lemma 7.1.8 and Proposition 7.1.10]. Let  $x_0 \in \mathbb{R}^{n+1}$  be a singular point with corresponding blow-up sequence  $(p_k, t_k)_{k \in \mathbb{N}}$  in  $M^n \times [0, T)$ . By the definition (9) of a type-I singularity, we calculate for  $p \in M^n$  and  $t_k, t_l \in [0, T)$ ,

$$|X(p,t_{l}) - X(p,t_{k})| \leq \int_{t_{k}}^{t_{l}} \left| \frac{\partial X}{\partial t}(p,t) \right| dt \leq \int_{t_{k}}^{t_{l}} |H(p,t)| dt$$

$$\leq 2 \int_{t_{k}}^{t_{l}} |H|_{\max}(t) dt \leq 2 \int_{t_{k}}^{t_{l}} \frac{C_{0}}{\sqrt{2(T-t)}} dt$$

$$= C_{0} \left( -\sqrt{2(T-t_{l})} + \sqrt{2(T-t_{k})} \right) \leq C_{0} \sqrt{2(T-t_{k})}. \quad (12)$$

Since the sequence  $(p_k)_{k\in\mathbb{N}}$  is bounded, there exist a point  $p_0\in M^n$  and a subsequence with

$$p_k \to p_0 \tag{13}$$

for  $k \to \infty$ . We employ (12) for  $p = p_l$ , and obtain

$$|X(p_l, t_l) - X(p_l, t_k)| \le C_0 \sqrt{2(T - t_k)}$$
 (14)

for all  $k, l \in \mathbb{N}$ . By Definition 5.1, we can choose  $l_0 = l_0(k)$  large enough so that, for fixed  $k \in \mathbb{N}$ ,

$$|X(p_l, t_l) - x_0| \le C_0 \sqrt{2(T - t_k)} \tag{15}$$

for all  $l \geq l_0$ . Estimates (14) and (15) imply

$$|X(p_l, t_k) - x_0| \le |X(p_l, t_k) - X(p_l, t_l)| + |X(p_l, t_l) - x_0|$$

$$\le 3C_0\sqrt{2(T - t_k)}$$
(16)

for fixed  $k \in \mathbb{N}$  and for all  $l \geq l_0(k)$ . For given  $\varepsilon > 0$ , choose  $k_0 = k_0(\varepsilon)$  large enough, so that

$$3C_0\sqrt{2(T-t_k)}<\frac{\varepsilon}{2}.$$

for all  $k \geq k_0$ . Then (16) yields

$$|X(p_l, t_k) - x_0| < \frac{\varepsilon}{2}$$

for all  $k \geq k_0(\varepsilon)$  and  $l \geq l_0(k)$ . By the convergence (13) and the continuity of the immersion X in its spatial argument, we can further choose  $l_0$  large enough, so that also

$$|X(p_0,t_k)-X(p_l,t_k)|<rac{arepsilon}{2}$$

for  $l \geq l_0$ . Hence,

$$|X(p_0, t_k) - x_0| \le |X(p_0, t_k) - X(p_{l_0}, t_k)| + |X(p_{l_0}, t_k) - x_0| < \varepsilon$$

for all  $k \geq k_0(\varepsilon)$ . Since  $\varepsilon > 0$  was chosen arbitrarily, we obtain

$$X(p_0, t_k) \to x_0 \tag{17}$$

for  $k \to \infty$ . Definition 5.1 and the type-I condition (9) yield

$$\lambda_k = |A(p_k, t_k)| \le \frac{C_0}{\sqrt{2(T - t_k)}}$$

and the estimate (12) implies

$$|X(p_0,t_l)-X(p_0,t_k)| \leq 2C_0\sqrt{2(T-t_k)} \leq \frac{2C_0^2}{\lambda_l}$$

We send  $l \to \infty$  in the above inequality and obtain with (17),

$$|\lambda_k|x_0 - X(p_0, t_k)| \le 2C_0^2$$

for all  $k \in \mathbb{N}$ . The definition (11) of the rescaled embedding provides, for  $\tau_k := \lambda_k^2(t_k - T)$ ,

$$|X_k(p_0, \tau_k)| = \lambda_k \left| X\left(p_0, T + \frac{\tau_k}{\lambda_k^2}\right) - x_0 \right| \le 2C_0^2$$

for all  $k \in \mathbb{N}$ . By the convergence (13), for given  $\delta > 0$ , there exists  $k_1 \in \mathbb{N}$  so that  $|p_k - p_0| < \delta$  for all  $k \ge k_0$ . By the continuity of the rescaled embedding, for given  $\varepsilon > 0$ , there exists  $\delta > 0$  so that, for  $|p_k - p_0| < \delta$ , we have

$$|X_k(p_k,\tau_k)-X_k(p_0,\tau_k)|<\varepsilon.$$

Hence, for given  $\varepsilon > 0$ , there exists  $k_1 \in \mathbb{N}$  so that

$$|X_k(0,\tau_k)| = |X_k(p_k,\tau_k)| \le |X_k(p_k,\tau_k) - X_k(p_0,\tau_k)| + |X_k(p_0,\tau_k)| < \varepsilon + 2C_0^2$$

for all  $k \geq k_1$ . Choosing  $\varepsilon = C_0^2$  yields  $X_k(0, \tau_k) \in B_{3C_0^2}(0)$  for all  $k \geq k_1$ . To bound the sequence

$$\left(\tau_k = -\lambda_k^2 (T - t_k)\right)_{k \in \mathbb{N}},$$

we estimate

$$\alpha_k = -\lambda_k^2 T < -\lambda_k^2 T + \lambda_k^2 t_k = \tau_k < 0$$

for all  $k \in \mathbb{N}$ . The rescaling behaviour from Remark 4.9 of the curvature yields

$$|A_k|^2(0,\tau_k) = |A_k|^2(p_k,\tau_k) = \frac{1}{\lambda_k^2}|A|^2\left(p_k,T + \frac{\tau_k}{\lambda_k^2}\right) = \frac{1}{\lambda_k^2}|A|^2(p_k,t_k) = 1.$$

Using Definition 5.1 and the lower blow-up rate from Proposition 4.6, we estimate

$$\tau_k = -\lambda_k^2 (T - t_k) = -|A|^2 (p_k, t_k) (T - t_k) \le -\frac{(T - t_k)}{2(T - t_k)} = -\frac{1}{2}$$

and, by the type-I assumption (9),

$$\tau_k = -\lambda_k^2 (T - t_k) = -|A|^2 (p_k, t_k) (T - t_k) \ge -\frac{C_0^2 (T - t_k)}{2(T - t_k)} = -\frac{C_0^2}{2}$$

for all  $k \in \mathbb{N}$ . For the curvature estimate, let  $\delta > 0$ ,  $k \in \mathbb{N}$ ,  $\tau \in [\alpha_k, -\delta^2]$  and  $p \in M^n$ . Then, the type-I condition (9) rescales to

$$|A_k(p,\tau)| = \frac{1}{\lambda_k} \left| A\left(p, T + \frac{\tau}{\lambda_k^2}\right) \right| \le \frac{1}{\lambda_k} \frac{C_0}{\sqrt{-2\tau/\lambda_k^2}} \le \frac{C_0}{\sqrt{-\tau}}.$$

Hence,

$$\max_{M^n \times [\alpha_k, -\delta^2]} |A_k| \le \frac{C_0}{\delta}$$

for each  $k \in \mathbb{N}$ .

**Theorem 5.3** (Convergence of rescalings). Let  $(M_t)_{t\in[0,T)}$  be a smooth, immersed solution of (MCF) with  $T < \infty$ . For the type-I rescaling 5.1 in case of a type-I singularity, there exists a sequence of rescaled immersions

$$\left( \left( M_{\tau}^{k} \right)_{\tau \in [\alpha_{k}, 0)} \right)_{k \in \mathbb{N}}$$

that converges for  $k \to \infty$  along a subsequence, uniformly and smoothly on compact subsets of  $(-\infty,0)$  and  $\mathbb{R}^{n+1}$  to a maximal, smooth limit solution  $(M_{\tau}^{\infty})_{\tau \in (-\infty,0)}$  which satisfies

$$M_{\tau_{\infty}}^{\infty} \cap B_{3C_0^2}(0) \neq \emptyset \qquad \text{ and } \qquad |A_{\infty}|^2(x) = 1 \ \text{ for some } \ x \in M_{\tau_{\infty}}^{\infty} \,,$$

where  $\tau_{\infty} \in [-C_0^2/2, -1/2]$  and, for  $\delta > 0$ ,

$$\sup_{\tau \in (-\infty, -\delta)} \sup_{M_\tau^\infty} |A_\infty| \le \frac{C_0}{\delta^2} \,.$$

Moreover, if  $(M_t)_{t\in[0,T)}$  is embedded, then  $(M_\tau^\infty)_{\tau\in(-\infty,0)}$  is embedded.

*Proof.* The convergence follows from Theorem 4.10 and Lemma 5.2 yields the properties. By Proposition 1.9,  $M_{\tau}^k$  is embedded for all  $k \in \mathbb{N}$  and all  $\tau \in [\alpha_k, 0)$ . Furthermore,

$$d_k(\tau) \ge \min \left\{ d_k(\alpha_k), \frac{\sin(\varepsilon)}{m_k(\tau)} \right\} \ge \min \left\{ \lambda_k d(0), \frac{\sin(\varepsilon)\delta^2}{C_0} \right\}$$

is uniformly bounded in k for  $\tau \leq \delta < 0$ .

5.1. Huisken's monotonicity formula. For  $x_0 \in \mathbb{R}^{n+1}$  and  $t_0 \in \mathbb{R}$ , define the backward heat kernel  $\Phi_{(x_0,t_0)} : \mathbb{R}^{n+1} \times (-\infty,t_0) \to \mathbb{R}$  by

$$\Phi_{(x_0,t_0)}(x,t) := \frac{1}{(4\pi(t_0-t))^{n/2}} \exp\left(-\frac{|x-x_0|^2}{4(t_0-t)}\right).$$

Let  $x, x_0, y_0 \in \mathbb{R}^{n+1}, t_0, \tau_0 \in \mathbb{R}, t \in (-\infty, t_0), \lambda > 0$  and  $\tau_0 > \lambda^2(t - t_0)$ . Then

$$\Phi_{(y_0,\tau_0)}(\lambda(x-x_0),\lambda^2(t-t_0)) = \frac{1}{\lambda^n} \Phi_{(x_0+y_0/\lambda,t_0+\tau_0/\lambda^2)}(x,t).$$

For the rescaled flow  $(M_{\tau}^{\lambda})_{\tau \in [-\lambda^2 T, 0)}$ ,

$$d\mu_{\lambda}^{n} = \sqrt{\det(g_{ij}^{\lambda})} dp = \sqrt{\lambda^{2n} \det(g_{ij})} dp = \lambda^{n} d\mu.$$

Hence, the integral

$$\int_{M_{\tau}^{\lambda}} \Phi_{(y_0, \tau_0)} d\mu_{\lambda}^n = \int_{M_{T-\tau/\lambda^2}} \Phi_{(x_0 + y_0/\lambda, T + \tau_0/\lambda^2)} d\mu^n$$

is scaling invariant, which makes it a useful quantity. In the following, we set H(x,t)=H(p,t) and  $\boldsymbol{\nu}(x,t)=\boldsymbol{\nu}(p,t)$  for x=X(p,t).

**Theorem 5.4** (Monotonicity formula, Huisken [Hui90, Theorem 3.1]). Let  $X: M^n \times (0,T) \to \mathbb{R}^{n+1}$  be a solution of (MCF). Then

$$\frac{d}{dt}\left(\int_{M_t}\Phi_{(x_0,t_0)}\,d\mu^n_t\right) = -\int_{M_t}\left|H - \frac{\langle x-x_0, \pmb{\nu}\rangle}{2(t_0-t)}\right|^2\Phi_{(x_0,t_0)}\,d\mu^n_t$$

for  $t_0 \in (0,T]$  and  $t \in (0,t_0)$ .

*Proof.* We follow the lines of [Hui90, Theorem 3.1]. We set  $x_0 = 0$  and  $t_0 = 0$ . Since x = x(t) with  $\partial_t x(t) = \mathbf{H}$ , we derive

$$\begin{split} \frac{d}{dt}\Phi_{(0,0)} &= \left(\frac{(n/2)4\pi}{-4\pi t} - \frac{2\langle x, \mathbf{H} \rangle}{-4t} - \frac{|x|^2}{4t^2}\right)\Phi_{(0,0)} \\ &= \left(\frac{n}{-2t} + H\frac{\langle x, \mathbf{\nu} \rangle}{-2t} - \frac{|x|^2}{4t^2}\right)\Phi_{(0,0)} \end{split}$$

so that

$$\frac{d}{dt} \left( \int_{M_t} \Phi_{(0,0)} d\mu_t^n \right) = \int_{M_t} \left( \frac{n}{-2t} + H \frac{\langle x, \boldsymbol{\nu} \rangle}{-2t} - \frac{|x|^2}{4t^2} - H^2 \right) \Phi_{(0,0)} d\mu_t^n$$

Observe that

$$-H^{2} + H\frac{\langle x, \boldsymbol{\nu} \rangle}{-t} - \frac{\langle x, \boldsymbol{\nu} \rangle^{2}}{4t^{2}} = -\left| H - \frac{\langle x, \boldsymbol{\nu} \rangle}{-2t} \right|^{2}$$

and

$$|x|^2 = \langle x, \boldsymbol{\nu} \rangle^2 + g^{ij} \langle x, \partial_i X \rangle \langle x, \partial_j X \rangle.$$

Hence,

$$\frac{n}{-2t} + H \frac{\langle x, \boldsymbol{\nu} \rangle}{-2t} - \frac{|x|^2}{4t^2} - H^2$$

$$= \frac{1}{-2t} \left( n - H \langle x, \boldsymbol{\nu} \rangle - \frac{1}{-2t} g^{ij} \langle x, \partial_i X \rangle \langle x, \partial_j X \rangle \right) - \left| H - \frac{\langle x, \boldsymbol{\nu} \rangle}{-2t} \right|^2 . \tag{18}$$

For  $x \in M_t$ ,

$$\operatorname{div}_{M_t} x = \operatorname{div}_{M^n} X(p, t) = n$$

and by the divergence theorem.

$$-\int_{M_t} H\langle x, \nu \rangle \Phi_{(0,0)} \, d\mu_t^n = \int_{M_t} \langle x, \mathbf{H} \rangle \Phi_{(0,0)} \, d\mu_t^n = -\int_{M_t} \operatorname{div}_{M_t} (x \, \Phi_{(0,0)}) \, d\mu_t^n \,,$$

where

$$\operatorname{div}_{M_t} \big( x \, \Phi_{(0,0)} \big) = \Phi_{(0,0)} \operatorname{div}_{M_t} x + \big\langle x, \nabla^{M_t} \Phi_{(0,0)} \big\rangle \; .$$

We calculate on  $M_t$ ,

$$\nabla^{M_t} \Phi_{(0,0)} = -\Phi_{(0,0)} g^{ij} \frac{2\langle x, \partial_i x \rangle}{-4t} \partial_j X = -\Phi_{(0,0)} g^{ij} \frac{\langle x, \partial_i X \rangle}{-2t} \partial_j X.$$

so that

$$\operatorname{div}_{M_t}(x \,\Phi_{(0,0)}) = n - \frac{1}{-2t} g^{ij} \langle x, \partial_i X \rangle \langle x, \partial_j X \rangle \Phi_{(0,0)}$$

which proves the claim.

**Theorem 5.5** (Weighted monotonicity formula, [Eck04, Theorem 4.13]). Let  $X: M^n \times (0,T) \to \mathbb{R}^{n+1}$  be a solution of (MCF) and  $\varphi: \mathbb{R}^{n+1} \times (0,T) \to \mathbb{R}$  in  $C^{2;1}$ . Then

$$\frac{d}{dt} \int_{M_t} \varphi \, \Phi_{(x_0, t_0)} \, d\mu_t^n = -\int_{M_t} \left| \mathbf{H} + \frac{(x - x_0)^{\perp}}{2(t_0 - t)} \right|^2 \varphi \, \Phi_{(x_0, t_0)} \, d\mu_t^n$$
$$+ \int_{M_t} \left( \frac{\partial}{\partial t} - \Delta_{M_t} \right) \varphi \, \Phi_{(x_0, t_0)} \, d\mu_t^n$$

for  $t_0 \in (0,T]$  and  $t \in (0,t_0)$ .

*Proof.* The proof is like the one for Theorem 5.4 with one additional step. When applying the divergence theorem, Theorem A.2, we now use the vector  $v = x\varphi \Phi_{(0,0)}$  instead and deduce

$$\int_{M_t} \langle x, H \boldsymbol{\nu} \rangle \varphi \, \Phi_{(0,0)} \, d\mu_t^n = \int_{M_t} \operatorname{div}_{M_t} \left( (x) \varphi \, \Phi_{(0,0)} \right) d\mu_t^n \,,$$

where

$$\operatorname{div}_{M_t}(x\varphi\,\Phi_{(0,0)}) = n\varphi\,\Phi_{(0,0)} + \varphi\,\langle x, \nabla^{M_t}\Phi_{(0,0)}\rangle + \langle x, \nabla^{M_t}\varphi\rangle\,\Phi_{(0,0)}.$$

Since  $\nabla^{M_t} \varphi = \tau_i(\varphi) \tau_i$  we can utilise the gradient of  $\Phi_{(0,0)}$  again to find

$$\frac{\left\langle x,\nabla^{M_t}\varphi\right\rangle}{-2t}\Phi_{(0,0)}=-\left\langle \nabla^{M_t}\Phi_{(0,0)},\nabla^{M_t}\varphi\right\rangle$$

so that integration by parts yields the extra term

$$\int_{M_t} \frac{\left\langle x, \nabla^{M_t} \varphi \right\rangle}{-2t} \Phi_{(0,0)} d\mu_t^n = \int_{M_t} \Delta_{M_t} \varphi \, \Phi_{(0,0)} d\mu_t^n.$$

The minus sign comes from the operation in (18).

**Remark 5.6** (see [Eck04, Remark 4.8]). If  $M_t$  is only defined locally, say in  $B_{\sqrt{4n\rho}}(x_0) \times (t_0 - \rho^2, t_0)$ , then we can use the cut-off function

$$\varphi_{(x_0,t_0)}^{\rho}(x,t) = \left(1 - \frac{|x - x_0|^2 + 2n(t_0 - t)}{\rho^2}\right)_{+}^{3}$$

where  $(\partial_t - \Delta_{M_t})\varphi \leq 0$ . Thus we still get the monotonicity inequality

$$\frac{d}{dt} \int_{M_t} \varphi_{(x_0,t_0)}^{\rho} \Phi_{(x_0,t_0)} d\mu_t^n \le -\int_{M_t} \left| \mathbf{H} + \frac{(x-x_0)^{\perp}}{2(t_0-t)} \right|^2 \varphi_{(x_0,t_0)}^{\rho} \Phi_{(x_0,t_0)} d\mu_t^n$$
 for  $t \in (0,t_0)$ .

**Theorem 5.7.** Let  $M_0$  be compact, convex and embedded. Then, every limit flow obtained by the type-I rescaling 5.1 around a type-I singularity, up to a rotation in  $\mathbb{R}^{n+1}$ , must be either the skrinking spheres  $(\mathbb{S}^n_{\sqrt{-2n\tau}})_{\tau \in (-\infty,0)}$  or one of the shrinking cylinders  $(\mathbb{S}^m_{\sqrt{-2m\tau}} \times \mathbb{R}^{n-m})_{\tau \in (-\infty,0)}$  for 0 < m < n.

*Proof.* Let  $x_0 \in \mathbb{R}^{n+1}$  be arbitrary. For  $t \in [0,T)$ , define the monotonicity quantity

$$\Theta_{(x_0,T)}(t) := \int_{M_t} \Phi_{(x_0,T)}(x,t) \, d\mu_t^n \, .$$

The monotonicity formula, Theorem 5.4, yields

$$\partial_t \Theta_{(x_0, T)}(t) \le 0 \tag{19}$$

for  $t \in (0,T)$ . Hence, the monotonicity quantity is monotonically decreasing and strictly positive, so that the limit

$$\lim_{t \to T} \Theta_{(x_0,T)}(t)$$

exists and for any sequence  $(t_k)_{k\in\mathbb{N}}$  with  $t_k\nearrow T$  for  $k\to\infty$ ,

$$\lim_{t \to T} \Theta_{(x_0,T)}(t) = \lim_{k \to \infty} \Theta_{(x_0,T)}(t_k). \tag{20}$$

For  $k \in \mathbb{N}$ ,  $y = \lambda_k(x - x_0) \in \mathbb{R}^{n+1}$  and  $\tau = \lambda_k^2(t - T) \in [\alpha_k, 0)$ , the backward heat kernel rescales according to

$$\Phi_{(0,0)}(y,\tau) = \frac{1}{\lambda_k^n} \Phi_{(x_0+0/\lambda_k,T+0/\lambda_k^2)}(x,t) = \frac{1}{\lambda_k^n} \Phi_{(x_0,T)}(x,t) \,.$$

Let  $\tau \in (-\infty, 0)$  and  $k_0 \in \mathbb{N}$  so that  $\tau \in [\alpha_k, 0)$  for  $k \geq k_0$ . Let  $(\lambda_k)_{k \in \mathbb{N}}$  be a sequence of positive real numbers with  $\lambda_k \to \infty$  for  $k \to \infty$ . We rescale the flow according to the type-I rescaling 5.1 with respect to the sequence  $(\lambda_k)_{k \in \mathbb{N}}$  and consider the rescaled flow  $(M_{\tau}^k)_{\tau \in [\alpha_k, 0)}$ . We receive a factor of  $\lambda_k^n$  from the scaling behaviour of the area element, and a factor of  $1/\lambda_k^n$  from the scaling behaviour of the backward heat kernel. Hence, the monotonicity quantity translates, for  $t_k := T + \tau/\lambda_k^2$ , by

$$\begin{split} \Theta_{(x_0,T)}(t_k) &= \int_{M_{t_k}} \Phi_{(x_0,T)}(x,t_k) \, d\mu_{t_k}^n \\ &= \int_{M_{\tau}^k} \Phi_{(0,0)}(y,\tau) \, d\mu_{k,\tau}^n =: \Theta_{(0,0)}^k(\tau) \, . \end{split}$$

Corollary 4.4 implies that there exist  $p_0 \in M^n$ ,  $x_0 \in \mathbb{R}^{n+1}$  and  $(p_k, t_k)_{k \in \mathbb{N}}$  with

$$X(p_k, t_k) \to x_0$$
 and  $|A(p_k, t_k)| = \max_{M_n} |A(\cdot, t_k)| \to \infty$ 

for  $k \to \infty$ . We rescale according to Definition 5.1 with respect to  $x_0$  and  $(p_k, t_k)_{k \in \mathbb{N}}$  and consider the rescaled embeddings  $X_k : M^n \times [\alpha_k, 0) \to \mathbb{R}^{n+1}$ . We apply the monotonicity formula 5.4 and estimate similar to [Bak10, Proposition 6.6] or [Coo11, Proposition 5.8],

$$0 \leq \int_{\tau_{1}}^{\tau_{2}} \int_{M_{\tau}^{k}} \left| \mathbf{H}_{k} + \frac{y^{\perp}}{-2\tau} \right|^{2} \Phi_{(0,0)} d\mu_{\tau}^{n} d\tau \leq \Theta_{(0,0)}^{k}(\tau_{1}) - \Theta_{(0,0)}^{k}(\tau_{2})$$

$$= \Theta_{(x_{0},T)} \left( T + \frac{\tau_{1}}{\lambda_{k}^{2}} \right) - \Theta_{(x_{0},T)} \left( T + \frac{\tau_{2}}{\lambda_{k}^{2}} \right)$$
(21)

for all  $k \geq k_0$ . Since

$$T + \frac{\tau_i}{\lambda_k^2} \to T$$

for  $k \to \infty$  and i = 1, 2, and by the existence of the limit (20), the right-hand side of (21) converges to 0 for  $k \to \infty$ . By Theorem 5.3, the sequence  $((M_{\tau}^k)_{\tau \in [\tau_1, \tau_2]})_{k \in \mathbb{N}}$  converges smoothly along a subsequence and on compact subsets of  $\mathbb{R}^{n+1}$  to a smooth flow  $(M_{\tau}^{\infty})_{\tau \in [\tau_1, \tau_2]}$ . Let R > 0. By the smooth convergence, there exists a  $k_0 \in \mathbb{N}$  so that for all  $k \geq k_0$ ,  $M_{\tau}^k \cap B_R(0)$  can be parametrized over  $M_{\tau}^{\infty} \cap B_R(0)$ . That is, there exist embeddings  $Y_k : M_{\tau}^{\infty} \cap B_R(0) \to \mathbb{R}^{n+1}$  with

$$M_{\tau}^k \cap B_R(0) = Y_k(M_{\tau}^{\infty} \cap B_R(0))$$

and  $Y_k \to \text{id}$  for  $k \to \infty$ . For  $\tau \in [\tau_1, \tau_2]$ , Fatou's lemma, Lemma ??, implies

$$0 = \liminf_{k \to \infty} \int_{M_{\tau}^{k} \cap B_{R}(0)} \left| \mathbf{H}_{k} + \frac{y^{\perp}}{-2\tau} \right|^{2} \Phi_{(0,0)} d\mu_{k,\tau}^{n}$$

$$= \liminf_{k \to \infty} \int_{M_{\tau}^{\infty} \cap B_{R}(0)} \left| \mathbf{H}_{k} + \frac{Y_{k}^{\perp}}{-2\tau} \right|^{2} \Phi_{(0,0)} \sqrt{\det(DY_{k})} dx$$

$$\geq \int_{M_{\tau}^{\infty} \cap B_{R}(0)} \liminf_{k \to \infty} \left( \left| \mathbf{H}_{k} + \frac{Y_{k}^{\perp}}{-2\tau} \right|^{2} \Phi_{(0,0)} \sqrt{\det(DY_{k})} \right) dx$$

$$= \int_{M_{\tau}^{\infty} \cap B_{R}(0)} \left| \mathbf{H}_{\infty} + \frac{Y_{\infty}^{\perp}}{-2\tau} \right|^{2} \Phi_{(0,0)} \sqrt{\det(DY_{\infty})} dx$$

$$= \int_{M_{\tau}^{\infty} \cap B_{R}(0)} \left| \mathbf{H}_{\infty} + \frac{y^{\perp}}{-2\tau} \right|^{2} \Phi_{(0,0)} d\mu_{\infty,\tau}^{n} \geq 0.$$

Thus also

$$\int_{\tau_1}^{\tau_2} \int_{M^{\infty} \cap B_R(0)} \left| \mathbf{H}_{\infty} + \frac{y^{\perp}}{-2\tau} \right|^2 \Phi_{(0,0)} d\mu_t^n d\tau = 0.$$

Since R > 0 was chosen arbitrarily, we deduce

$$\int_{\tau_1}^{\tau_2} \int_{M^{\infty}} \left| \mathbf{H}_{\infty} + \frac{y^{\perp}}{-2\tau} \right|^2 \Phi_{(0,0)} d\mu_t^n d\tau = 0.$$

Since the convergence is smooth, and sending  $\tau_1 \to -\infty$  and  $\tau_2 \to 0$  yields

$$\left| \mathbf{H}_{\infty} + \frac{y^{\perp}}{-2\tau} \right|^2 = 0$$

for every  $\tau \in (-\infty, 0)$  and every  $y \in M_{\tau}^{\infty}$ .

For the area estimate, let again be R > 0 and  $\tau \in (-\infty, 0)$ . Then there exists again  $k_0 \in \mathbb{N}$  so that  $\tau \in [\alpha_k, 0)$  and

$$T - \frac{\tau}{\lambda_k^2} \ge \frac{T}{2}$$

for all  $k \geq k_0$ . Like in Corollary 1.5,

$$\partial_t \mu_t^n(M_t \cap B_R) = -\int_{M_t \cap B_R} H^2 d\mu_t^n \,,$$

the area is decreasing locally also locally. By (19), the monotonicity quantity is decreasing in time and we can estimate with the definition of the backward heat kernel and the behaviour of the area of the hypersurfaces,

$$\begin{split} & \int_{M_{\tau}^{k} \cap B_{R}(0)} \Phi_{(0,0)}(y,\tau) \, d\mu_{k,\tau}^{n} \\ & \leq \int_{M_{T-\tau/\lambda_{k}^{2}} \cap B_{R}(x_{0})} \Phi_{(x_{0},T)} \bigg( x,T - \frac{\tau}{\lambda_{k}^{2}} \bigg) \, d\mu_{T-\tau/\lambda_{k}^{2}}^{n} \\ & \leq \int_{M_{T/2} \cap B_{R}(x_{0})} \Phi_{(x_{0},T)} \bigg( x,\frac{T}{2} \bigg) \, d\mu_{T/2}^{n} \\ & = \frac{1}{(4\pi(T-T/2))^{n/2}} \int_{M_{T/2} \cap B_{R}(x_{0})} \exp \bigg( -\frac{|x-x_{0}|^{2}}{4(T-T/2)} \bigg) \, d\mu_{T/2}^{n} \\ & \leq C(n,T) \mu_{T/2}^{n}(M_{T/2} \cap B_{R}(x_{0})) \leq C(n,T) \mu_{0}^{n}(M_{0} \cap B_{R}(x_{0})) \end{split}$$

Like before, Fatou's lemma implies

$$\liminf_{k \to \infty} \int_{M_{\tau}^{k} \cap B_{R}(0)} \Phi_{(0,0)} d\mu_{k,\tau}^{n}$$

$$= \liminf_{k \to \infty} \int_{M_{\tau}^{\infty} \cap B_{R}(0)} \Phi_{(0,0)} \sqrt{\det(DY_{k})} dx$$

$$\geq \int_{M_{\tau}^{\infty} \cap B_{R}(0)} \liminf_{k \to \infty} \left( \Phi_{(0,0)} \sqrt{\det(DY_{k})} \right) dx$$

$$= \int_{M_{\infty}^{\infty} \cap B_{R}(0)} \Phi_{(0,0)} d\mu_{\infty,\tau}^{n} .$$

Furthermore,

$$\begin{split} & \int_{M_{\tau}^{\infty} \cap B_{R}(0)} \Phi_{(0,0)}(y,\tau) \, d\mu_{\infty,\tau}^{n} \\ & = \frac{1}{(-4\pi\tau)^{n/2}} \int_{M_{\tau}^{\infty} \cap B_{R}(0)} \exp\left(-\frac{|y|^{2}}{-4\tau}\right) d\mu_{\infty,\tau}^{n} \\ & \geq \frac{1}{(-4\pi\tau)^{n/2}} \int_{M_{\tau}^{\infty} \cap B_{R}(0)} \exp\left(-\frac{R^{2}}{-4\tau}\right) d\mu_{\infty,\tau}^{n} \\ & = \frac{1}{(-4\pi\tau)^{n/2}} \exp\left(-\frac{R^{2}}{-4\tau}\right) \mu_{\infty,\tau}^{n}(M_{\tau}^{\infty} \cap B_{R}(0)) \, . \end{split}$$

so that

$$\mu^n(M_{\tau}^{\infty} \cap B_R(0)) \le C(n, T, \tau)\mu_0^n(M_0 \cap B_R(x_0)) \exp\left(\frac{R^2}{-4\tau}\right)$$

holds for all  $\tau \in (-\infty, 0)$ .

For every fixed  $\tau \in (-\infty, 0)$ , by Theorem 5.3, |A| is not identically zero and  $|\nabla^m A| \leq C_m$ , for every  $m \in \mathbb{N}$ . Theorem 2.4 yields that

$$M_{\tau}^{\infty} = \mathbb{S}_{\sqrt{-2m\tau}}^{m} \times \mathbb{R}^{n-m}$$

where  $0 < m \le n$ . Since the flow is smooth, the claim follows.

# 6. Typ-II singularities

The following rescaling technique for type-II singularities was introduced in [Ham95a, Proof of Theorem 16.4] for Ricci flow, and applied to type-II singularities of MCF in [HS99, p. 11].

**Definition 6.1** (Type-II rescaling). Let  $(p_k, t_k)_{k \in \mathbb{N}}$  be a sequence in  $M^n \times [0, T - 1/k]$  with

$$H^{2}(p_{k}, t_{k}) \left( T - \frac{1}{k} - t_{k} \right) = \max_{(p, t) \in M^{n} \times [0, T - 1/k]} \left( H^{2}(p, t) \left( T - \frac{1}{k} - t \right) \right)$$

for each  $k \in \mathbb{N}$ . We set

$$\lambda_k^2 := |A|^2(p_k, t_k), \qquad \alpha_k := -\lambda_k^2 t_k \qquad \text{and} \qquad T_k := \lambda_k^2 \left( T - \frac{1}{k} - t_k \right).$$

and define the rescaled embeddings  $X_k: M^n \times [\alpha_k, T_k] \to \mathbb{R}^2$  by

$$X_k(p,\tau) := \lambda_k \left( X \left( p, t_k + \frac{\tau}{\lambda_k^2} \right) - X(p_k, t_k) \right) .$$

**Lemma 6.2** (Properties of the type-II rescaling, [HS99, Lemma 4.3]). Let  $X : M^n \times (0,T) \to \mathbb{R}^2$  be a solution of (MCF) with  $T < \infty$ . For the type-II rescaling 6.1 in case of a type-II singularity,

$$\lambda_k \to \infty$$
,  $\alpha_k \to -\infty$  and  $T_k \to \infty$ 

for  $k \to \infty$ . Moreover,

$$X_k(0,0) = 0$$
 and  $|A_k|^2(0,0) = 1$ 

for every  $k \in \mathbb{N}$  and for any  $\varepsilon > 0$  and any  $\overline{T} > 0$ , there exists a  $k_0 \in \mathbb{N}$  such that

$$\max_{M^n \times [\alpha_k, \bar{T}]} |A_k|^2 < 1 + \varepsilon$$

for all  $k \geq k_0$ .

*Proof.* We follow the lines of [HS99, Lemma 4.3]. By definition,  $X_k(0,0) = X_k(p_k,0) = 0$  and

$$|A_k|^2(0,0) = \frac{1}{\lambda_k^2}|A|^2(p_k, t_k) = 1$$

for each  $k \in \mathbb{N}$ . Let m > 0 be arbitrary. By the definition (10) of a type-II singularity, there exist  $\bar{t} \in [0, T)$  and  $\bar{p} \in M^n$  so that

$$|A|^2(\bar{p},\bar{t})(T-\bar{t}) > 2m$$
.

We fix  $\bar{t}$  and choose  $k_0 \in \mathbb{N}$ , so that  $\bar{t} < T - 1/k$  and  $|A|^2(\bar{p}, \bar{t})/k < m$  for all  $k \ge k_0$ . Then

$$|A|^2(\bar{p},\bar{t})\left(T-\frac{1}{k}-\bar{t}\right) = |A|^2(\bar{p},\bar{t})(T-\bar{t}) - \frac{1}{k}|A|^2(\bar{p},\bar{t}) > m$$

and Definition 6.1 yields

$$T_k = |A|^2(p_k, t_k) \left(T - \frac{1}{k} - t_k\right) \ge |A|^2(\bar{p}, \bar{t}) \left(T - \frac{1}{k} - \bar{t}\right) > m.$$

Since m was chosen arbitrarily, it follows that  $T_k \to \infty$  and thus also  $\lambda_k = |A|^2(p_k, t_k) \to \infty$  for  $k \to \infty$ . Since  $t_k \nearrow T$ , we conclude that  $\alpha_k = -\lambda_k^2 t_k \to -\infty$  for  $k \to \infty$ . For the curvature estimate, it again follows from Definition 6.1 that

$$|A|^2(p,t)\left(T - \frac{1}{k} - t\right) \le |A|^2(p_k, t_k)\left(T - \frac{1}{k} - t_k\right) = T_k$$
 (22)

for all  $p \in M^n$ ,  $t \in [0, T - 1/k]$  and  $k \in \mathbb{N}$ . Let  $\varepsilon > 0$  and  $\overline{T} > 0$  be given. Since  $T_k \to \infty$ , there exists again  $k_1 \in \mathbb{N}$  so that, for all  $k \geq k_1$ ,  $\overline{T} < T_k$  and

$$0 < \frac{\bar{T}}{T_k - \bar{T}} < \varepsilon.$$

For  $\tau \in [\alpha_k, \bar{T}]$ , it is  $t := t_k + \tau/\lambda_k^2 \in [0, T - 1/k)$ , and we can use the scaling behaviour of the curvature and (22) to estimate

$$|A_k|^2(p,\tau) = \frac{1}{\lambda_k^2} |A|^2 \left( p, t_k + \frac{\tau}{\lambda_k^2} \right) \le \frac{T - 1/k - t_k}{T - 1/k - (t_k + \tau/\lambda_k^2)}$$
$$= \frac{T_k}{T_k - \tau} \le \frac{T_k}{T_k - \bar{T}} = 1 + \frac{\bar{T}}{T_k - \bar{T}} < 1 + \varepsilon$$

for all  $p \in M^n$  and  $k \ge k_1$ . Hence,

$$\max_{M^n \times [\alpha_k, \bar{T}]} |A_k|^2 < 1 + \varepsilon$$

for all  $k \ge \max\{k_0, k_1\}$ .

**Theorem 6.3.** Let  $(M_t)_{t \in [0,T)}$  be a smooth, immersed solution of (MCF) with  $T < \infty$ . For the type-II rescaling 6.1 in case of a type-II singularity, there exists a sequence of rescaled immersions

$$\left( \left( M_{\tau}^{k} \right)_{\tau \in [\alpha_{k}, T_{k}]} \right)_{k \in \mathbb{N}}$$

that converges for  $k \to \infty$  along a subsequence, uniformly and smoothly on compact subsets of  $\mathbb{R}$  and  $\mathbb{R}^{n+1}$  to a maximal, smooth limit solution  $(M_{\tau}^{\infty})_{\tau \in \mathbb{R}}$  which satisfies again (MCF) and

$$0 \in M_0^{\infty}$$
 and  $\sup_{\mathbb{R} \times \mathbb{R}} |A_{\infty}| = |A_{\infty}(0)| = 1$ .

Moreover, if  $(M_t)_{t\in[0,T)}$  is embedded, then  $(M_{\tau}^{\infty})_{\tau\in(-\infty,0)}$  is embedded.

*Proof.* The convergence follows from Theorem 4.10. Lemma 6.2 implies  $0 \in M_0^{\infty}$  and  $|A_{\infty}(0)| = 1$  and that for any  $\varepsilon > 0$  and any  $\bar{T} > 0$ ,

$$\sup_{\mathbb{R}\times(-\infty,\bar{T}]}|A_{\infty}|^2\leq 1+\varepsilon.$$

Sending  $\bar{T} \to \infty$  and  $\varepsilon \to 0$  yields

$$\sup_{\mathbb{R}\times\mathbb{R}}|A_{\infty}|\leq 1=|A_{\infty}(0)|.$$

By Proposition 1.9,  $M_{\tau}^k$  is embedded for all  $k \in \mathbb{N}$  and all  $\tau \in [\alpha_k, T_k]$ . Furthermore,

$$d_k(\tau) \ge \min \left\{ d_k(\alpha_k), \frac{\sin(\varepsilon)}{m_k(\tau)} \right\} \ge \min \{ \lambda_k d(0), \sin(\varepsilon) \}$$

is uniformly bounded in k for  $\tau \in \mathbb{R}$ .

**Remark 6.4.** In the following chapters, we will show that the eternal solution obtained in Theorem 6.3 is convex and translating.

## 7. Convex hypersurfaces

**Theorem 7.1** (Huisken, [Hui84, Corollary 4.2]). Assume  $M_0 = X_0(M)$  closed and convex, i.e.  $h_{ij} \succeq 0$ . Then  $h_{ij} \succ 0$  for all  $t \in (0,T)$ .

Proof. By Lemma 1.4 and Simons' identity

$$\partial_t h_{ij} = \Delta h_{ij} - 2Hg^{km}h_{ik}h_{jm} + |A|^2 h_{ij}.$$

Use Theorem C.5 for  $m_{ij} = h_{ij}$ ,  $u^k \equiv 0$  and  $b_{ij} = -2Hh_{il}g^{lm}h_{mj} + |A|^2h_{ij}$ .

**Corollary 7.2.** There is some  $\varepsilon > 0$  such that  $h_{ij} \succeq \varepsilon Hg_{ij}$  holds on  $M \times (0,T)$ .

**Theorem 7.3** (Huisken, [Hui84, Theorem 4.3]). If  $\varepsilon H g_{ij} \leq h_{ij} \leq \beta H g_{ij}$ , and H > 0 at t = 0 for some constants  $0 < \varepsilon \leq 1/n < \beta < 1$ , then this remains so on (0,T).

*Proof.* To prove the first inequality, we want to apply Theorem C.5 with

$$m_{ij} = \frac{h_{ij}}{H} - \varepsilon g_{ij} , \quad u^k = \frac{2}{H} g^{kl} \nabla_l H , \quad b_{ij} = 2\varepsilon H h_{ij} - 2h_{im} g^{ml} h_{lj} .$$

With this choice the evolution equation in Theorem C.5 is satisfied since

$$\partial_t \left( \frac{h_{ij}}{H} \right) = \frac{1}{H^2} \left( H \Delta h_{ij} - h_{ij} \Delta h_{ij} \right) - 2h_{im} g^{ml} h_{mj}$$

and

$$\Delta \left(\frac{h_{ij}}{H}\right) = \frac{1}{H^2} \left( H \Delta h_{ij} - h_{ij} \Delta h_{ij} \right) - \frac{2}{H} g^{kl} \nabla_k H \nabla_l \left(\frac{h_{ij}}{H}\right) \,.$$

It remains to check that  $b_{ij}$  is nonnegative on the null-eigenvectors of  $m_{ij}$ . Assume that, for some vector v,

$$h_{ij}v^j = \varepsilon H v_i$$
.

Then we derive

$$b_{ij}v^{i}v^{j} = 2\varepsilon H h_{ij}v^{i}v^{j} - 2h_{im}g^{mI}h_{ij}v^{i}v^{j} = 2\varepsilon^{2}H^{2}|v|^{2} - 2\varepsilon^{2}H^{2}|v|^{2} = 0.$$

That the second inequality remains true follows in the same way after reversing signs.  $\Box$ 

**Theorem 7.4** (Huisken [Hui84]). Let  $n \geq 2$  and  $M_0 \subset \mathbb{R}^{n+1}$  be closed, convex and embedded. Then the mean curvature flow  $(M_t)_{t \in [0,T)}$  starting at  $M_0$  converges to a round point.

Proof. See [Man11, Theorem 3.4.10]. Let T be the maximal time of smooth existence of the mean curvature flow of an n-dimensional convex hypersurface. By Theorems 1.10, 7.1 and 7.3, we have that after any positive time H > 0 and there exists  $\varepsilon > 0$ , independent of time, such that  $h_{ij} \succeq \varepsilon H g_{ij}$ . If at time T we have a type-II singularity, we get an unbounded, eternal convex blow-up limit flow with  $H \geq 0$ , using Hamiltons procedure. By the strong maximum principle, actually H > 0 for every time (otherwise  $H \equiv 0$ , but this and the convexity would imply that the limit flow is simply a fixed hyperplane) and the condition  $h_{ij} \succeq \varepsilon H g_{ij}$  passes to the limit. Then, by Theorem 3.6, all the hypersurfaces of the limit flow are compact, in contradiction with the unboundedness, hence type-II singularities cannot develop. Dealing with type-I singularities, any blow-up limit is embedded, strictly convex and compact, again by this theorem. Hence, by Theorem 5.7 it can be only the sphere  $\mathbb{S}^n$ . This implies that the full sequence of rescaled hypersurfaces converges in  $C^{\infty}$  to such sphere. Finally, as the blow-up limit is unique and compact, the original hypersurface shrinks to a point in finite time.

Remark 7.5 (Exponential convergence, [Hui84, Lemma 10.6]). Consider the normalized flow

$$\tilde{X}(\cdot,t)=\psi(t)X(\cdot,t)$$

where  $\psi$  is chosen so that

$$\int_{\tilde{M}_t} d\tilde{\mu} = |M_0|$$

for all  $t \in [0,T)$ . By choosing

$$\tilde{t}(t) = \int_0^t \psi^2(\tau) \, d\tau \,,$$

we get that  $\tilde{g}_{ij} = \psi^2 g_{ij}$ ,  $\tilde{H} = \psi^{-1} H$ ,

$$\psi^{-1}\partial_t \psi = \frac{\int_{\tilde{M}_t} H^2 d\tilde{\mu}}{n \int_{\tilde{M}_t} d\tilde{\mu}} =: \frac{h}{n} = \psi^{-2} \frac{\tilde{h}}{n}$$

and

$$\partial_{\tilde{t}}\tilde{X} = \psi^{-2}\partial_{t}\tilde{X} = -\tilde{H}\tilde{\nu} + \frac{\tilde{h}}{n}\tilde{X}$$

for  $\tilde{t} \in [0, \infty)$ . Then there exist constants  $\delta > 0$  and  $C, C_m < \infty$  such that

$$\begin{aligned} \tilde{H}_{\max} - \tilde{H}_{\min} &\leq C e^{-\delta \tilde{t}} ,\\ \left| \tilde{h}_{ij} \tilde{H} - \frac{\tilde{h}}{n} \tilde{g}_{ij} \right| &\leq C e^{-\delta \tilde{t}} ,\\ \max_{\tilde{M}} \left| \nabla^m \tilde{A} \right| &\leq C_m e^{-\delta \tilde{t}} \end{aligned}$$

for all m > 0.

### 8. Hamilton's Harnack Inequality

We follow [Urb91, Section 2], [And94] and [Sch17c, Chapter 4]. For convex hypersurfaces, the initial value problem (MCF) can be reduced to an initial value problem for the support function. Let M be a smooth, closed, stricly convex hypersurface (A is positive definite everywhere). Recall the Gauss map  $\boldsymbol{\nu}: M^n \to \mathbb{S}^n$ , unit normal  $\bar{\boldsymbol{\nu}}: M^n \to \mathbb{R}^{n+1}$  and the Weingarten map  $S: TM^n \to TM^n$  which gives the rate of change in the direction of the normal along the surface with

$$S(v) := dX^{-1} (D_{dX(v)} \bar{\boldsymbol{\nu}}) = dX^{-1} (d_v \boldsymbol{\nu}).$$

The second fundamental form A is the symmetric tensor given by the normal component of the connection on  $\mathbb{R}^{n+1}$ .

$$A(u,v) = -\langle d^2X(v,w), \bar{\boldsymbol{\nu}}\rangle = -\langle D_{dX(v)}dX(w), \bar{\boldsymbol{\nu}}\rangle$$
$$= \langle dX(w), D_{dX(v)}\bar{\boldsymbol{\nu}}\rangle = g(w,S(v))$$

for all  $v,w\in TM^n$ , where  $dX:TM^n\to\mathbb{R}^{n+1}$ . The eigenvalues  $\lambda_1\ldots\lambda_n$  of S are called the principal curvatures. Without loss of generality, we may assume that M encloses the origin. All information about the hypersurface is contained in the support function  $s:M^n\to\mathbb{R}$  where

$$s(p) := \langle \bar{\boldsymbol{\nu}}(p), X(p) \rangle$$
.

For strictly convex hypersurfaces  $\nu$  is a global diffeomorphism, and we can parametrise the hypersurface by  $\tilde{X} : \nu(M^n) \subset \mathbb{S}^n \to \mathbb{R}^{n+1}$  where

$$\tilde{X}(z) := X(\boldsymbol{\nu}^{-1}(z))$$

for all  $z \in \nu(M^n)$ . We will consider the support function

$$s(z) := \langle \bar{z}, \tilde{X}(z) \rangle. \tag{23}$$

In the following, Indetify  $\bar{z}$  with z. If the support function is known, the hypersurface is given as the boundary of the convex region

$$\bigcap_{z \in \mathbb{S}^n} \left\{ y \in \mathbb{R}^{n+1} \, | \, \langle y, z \rangle \le s(z) \right\} \, .$$

Let  $\sigma_{ij}$  be the metric and  $\tilde{\nabla}$  be the gradient on  $\mathbb{S}^n$ . Differentiating (23) we obtain

$$\tilde{\nabla}_{i}s = \langle \tilde{\nabla}_{i}\tilde{X}, z \rangle + \langle \tilde{X}, \tilde{\nabla}_{i}z \rangle = \langle \tilde{X}, \tilde{\nabla}_{i}z \rangle,$$

since  $\tilde{\nabla}_i \tilde{X}(z)$  is tangential to M at  $\tilde{X}(z)$ , and z is the normal to M at  $\tilde{X}(z)$ . Since  $\langle z, z \rangle = 1$ , we obtain

$$\langle z, \tilde{\nabla}_i z \rangle = 0$$

and writing  $\tilde{\nabla}_{ij} := \tilde{\nabla}_i \tilde{\nabla}_j$ , we obtain

$$\langle z, \tilde{\nabla}_{ij} z \rangle = -\langle \tilde{\nabla}_{j} z, \tilde{\nabla}_{i} z \rangle = -\sigma_{ij}$$

Hence,

$$\begin{split} \tilde{X} &= \langle \tilde{X}, z \rangle z + \sigma^{ij} \langle \tilde{X}, \tilde{\nabla}_i z \rangle \tilde{\nabla}_j z \\ &= sz + \sigma^{ij} \tilde{\nabla}_i s \nabla_j z = sz + \tilde{\nabla} s \end{split}$$

From this, we conclude at a fixed point

$$\begin{split} \tilde{\nabla}_{i}\tilde{X} &= \tilde{\nabla}_{i}sz + s\tilde{\nabla}_{i}z + \tilde{\nabla}_{ki}s\sigma^{kl}\tilde{\nabla}_{l}z + \tilde{\nabla}_{k}s\sigma^{kl}\tilde{\nabla}_{li}z \\ &= \tilde{\nabla}_{i}sz + s\tilde{\nabla}_{i}z + \tilde{\nabla}_{ki}s\sigma^{kl}\tilde{\nabla}_{l}z - \tilde{\nabla}_{k}s\sigma^{kl}\sigma_{li}z \\ &= s\tilde{\nabla}_{i}z + \tilde{\nabla}_{ki}s\sigma^{kl}\tilde{\nabla}_{l}z \end{split}$$

and

$$\tilde{\nabla}_{ij}\tilde{X} = \tilde{\nabla}_{j}s\tilde{\nabla}_{i}z - s\sigma_{ij}z + \tilde{\nabla}_{kij}s\sigma^{kl}\tilde{\nabla}_{l}z - \tilde{\nabla}_{ki}s\sigma^{kl}\sigma_{lj}z$$
$$= \tilde{\nabla}_{j}s\tilde{\nabla}_{i}z - s\sigma_{ij}z + \tilde{\nabla}_{kij}s\sigma^{kl}\tilde{\nabla}_{l}z - \tilde{\nabla}_{ij}sz$$

so that

$$\tilde{h}_{ij} = -\langle \tilde{\nabla}_{ij} \tilde{X}, z \rangle = s \sigma_{ij} + \tilde{\nabla}_{ij} s$$

and

$$\tilde{g}_{ij} = s^2 \sigma_{ij} + 2s \tilde{\nabla}_{ij} s + \tilde{\nabla}_{ik} s \sigma^{kl} \tilde{\nabla}_{jl} s = \tilde{h}_{ik} \sigma^{kl} \tilde{h}_{lj}$$

as well as

$$\tilde{h}_{i}^{j} = \tilde{g}^{jk} \tilde{h}_{ik} = \tilde{a}^{jl} \sigma_{lm} \tilde{a}^{mk} \tilde{h}_{ik} = \sigma_{il} \tilde{a}^{lj}$$

where here  $(\tilde{a}^{ij})_{ij} = ((\tilde{h}_{ij})_{ij})^{-1}$  and

$$\tilde{H} = \tilde{h}_i^i = \sigma_{ij}\tilde{a}^{ij}$$

We consider the Weingarten map  $\tilde{S}: T\mathbb{S}^n \to T\mathbb{S}^n$  with

$$\tilde{S}(v) := d\tilde{X}^{-1}(d_v\tilde{\boldsymbol{\nu}}).$$

Since  $d\tilde{\nu} = \mathrm{id}$ , we have  $\tilde{S}^{-1} = d\tilde{X}$ . We define

$$\tilde{S}^{-1}(v) = (\sigma^* \tilde{\nabla}^2 s + s \operatorname{id})(v) = \tilde{\nabla}_v(\tilde{\nabla} s) + s \operatorname{id}(v) =: \mathcal{A}(v)$$
(24)

so that

$$\tilde{g}(u,v) = \tilde{g}_{ij}v^iw^j = \tilde{h}_{ik}\sigma^{kl}\tilde{h}_{lj}v^iw^j = \sigma_{km}\tilde{a}_i^m\sigma^{kl}\sigma_{ln}\tilde{a}_j^nv^iw^j$$
$$= \sigma_{km}\tilde{a}_i^m\tilde{a}_i^kv^iw^j = \sigma(\mathcal{A}(u),\mathcal{A}(v)).$$

The great advantage of the support function is that it allows us to consider a family of convex hypersurfaces simply as an evolving scalar function defined on the sphere. This makes things much simpler than the more abstract framework allowing arbitrary parametrizations, since we no longer have different descriptions of the same hypersurface. Furthermore, the identification with the sphere provides a time-independent metric and connection, which vastly simplifies many calculations, including especially those presented here for the proof of the Harnack inequalities.

For the remainder of this section, we consider a familiy of embeddigns  $X: M^n \times [0,T) \to \mathbb{R}^{n+1}$  that solve the initial value problem

$$\begin{cases} \partial_t X(p,t) = -F(S(p,t), \boldsymbol{\nu}(p,t)) \boldsymbol{\nu}(p,t) & \text{for } (p,t) \in M^n \times [0,T) \\ X(\cdot,0) = X_0 & \text{on } M^n \,. \end{cases}$$
 (25)

where F is such that the equation is parabolic and invariant under diffeomorphisms of  $M^n$  and translations in space and time. We want to reduce (25) to an initial value problem for the support function. Let X be a solution of (25), and suppose that for each  $t \in [0,T)$ ,  $X(\cdot,t)$  is a parametrization of a smooth, closed, uniformly convex hypersurface  $M_t$ . We define a new parametrization  $\tilde{X}(\cdot,t)$  by

$$\tilde{X}(z,t) = X\big(\boldsymbol{\nu}_t^{-1}(z),t\big) \ .$$

Then

$$\partial_t \tilde{X} = \partial_i X \partial_t (\boldsymbol{\nu}_t^{-1})^i + \partial_t X = \partial_i X \partial_t (\boldsymbol{\nu}_t^{-1})^i - \tilde{F}z$$

so that

$$\partial_t s = \langle \partial_t \tilde{X}, z \rangle = -\tilde{F}$$

since  $\partial_i X$  is tangential. This proves the following theorem:

**Theorem 8.1** (Andrews, [And94, Theorem 3.1]). Suppose  $X: M^n \times [0,T) \to \mathbb{R}^{n+1}$  is a family of strictly convex immersions satisfying (25). Then

$$\begin{cases} \partial_t s(z,t) = \Phi(\mathcal{A}[s(z,t)], z) & on \, \mathbb{S}^n \times [0,T) \\ s(\cdot,0) = s_0 & on \, \mathbb{S} \,. \end{cases}$$
 (26)

where id is the identity matrix,  $s_0$  is the support function of  $M_0$ ,

$$\Phi(\mathcal{A}) = -\operatorname{tr}_{\sigma} \mathcal{A}^{-1}$$
 and  $\mathcal{A} = \sigma^* \tilde{\nabla}^2 s + \operatorname{id} s$ .

The expression (24) allows us to use the support function to calculate functions of the curvature of a hypersurface. We can define  $\Phi: U \subset T^*\mathbb{S}^n \to \mathbb{R}$  in terms of X by

$$\Phi(X) = -\tilde{F}(\tilde{X}^{-1})$$

for all positive definite maps X. Furthermore,  $\dot{\Phi}(A): T^*\mathbb{S}^n \to T\mathbb{S}^n$  is given by

$$\dot{\Phi}(\mathcal{A})(\mathcal{B}) = \partial_r|_{r=0} \Phi(\mathcal{A} + r\mathcal{B})$$

and  $\ddot{\Phi}(\mathcal{A}): T\mathbb{S}^n \otimes T^*\mathbb{S}^n \to T\mathbb{S}^n \otimes T^*\mathbb{S}^n$  by

$$\ddot{\Phi}(\mathcal{A})(\mathcal{B},\mathcal{C}) = \partial_r|_{r=0}\dot{\Phi}(\mathcal{A} + r\mathcal{C})(\mathcal{B}).$$

We call  $\Phi$  concave (convex), if

$$\ddot{\Phi}(\mathcal{A})(\mathcal{B},\mathcal{B}) \leq (\geq) 0$$

for all  $\mathcal{A}, \mathcal{B} \in T^*\mathbb{S}^n$ . We call  $\Phi$   $\alpha$ -concave ( $\alpha$ -convex), if

$$\Phi = \operatorname{sign} \alpha B^{\alpha},$$

where B is positive and concave (convex),  $\alpha \in \mathbb{R}$ .  $\alpha$ -concavity (-convexity) is equivalent to

$$\ddot{\Phi} = \alpha(\alpha - 1)B^{\alpha - 2}\dot{B}\otimes\dot{B} + \alpha B^{\alpha - 1}\ddot{B} \preceq (\succeq)\frac{\alpha - 1}{\alpha\Phi}\dot{\Phi}\otimes\dot{\Phi}.$$
 (27)

(These conditions become considerably more complicated when written in terms of the principal curvatures and a speed function F. For example, concavity of  $\Phi$ , becomes  $\ddot{F}(X,X) + 2\dot{F}(X \circ S^{-1} \circ X) \geq 0$ .)

**Lemma 8.2** (Andrews, [And94, Theorem 3.6 and Lemma 5.1]). The following evolution equations hold under the Gauss map parametrization of the flow (25):

$$\partial_{t}(\tilde{\nabla}^{2}s + s\sigma) = \tilde{\nabla}^{2}\Phi + \Phi\sigma$$

$$\partial_{t}\mathcal{A} = \sigma^{*}\tilde{\nabla}^{2}\Phi + \Phi \operatorname{id}$$

$$\partial_{t}\Phi(\mathcal{A}) = \dot{\Phi}(\mathcal{A})(\sigma^{*}\tilde{\nabla}^{2}\Phi) + \dot{\Phi}(\mathcal{A})(\operatorname{id})\Phi$$

$$\partial_{t}^{2}\Phi(\mathcal{A}) = \ddot{\Phi}(\mathcal{A})(\partial_{t}\mathcal{A}, \partial_{t}\mathcal{A}) + \dot{\Phi}(\mathcal{A})(\sigma^{*}\tilde{\nabla}^{2}\partial_{t}\Phi) + \dot{\Phi}(\mathcal{A})(\operatorname{id})\partial_{t}\Phi .$$
(28)

*Proof.* The first equation follows simply by differentiating (26), since the metric  $\sigma$  and connection  $\tilde{\nabla}$  are independent of time. The second follows immediately from this. Since  $\Phi$  depends only on  $\mathcal{A}$ , we have  $\partial_t \Phi = \dot{\Phi}(\partial_t \mathcal{A})$  which implies the third equation. By (28),

$$\begin{split} \partial_t^2 \Phi &= \partial_t \Big( \dot{\Phi}(\sigma^* \tilde{\nabla}^2 \Phi) + \dot{\Phi}(\mathrm{id}) \Phi \Big) \\ &= \ddot{\Phi}(\partial_t \mathcal{A}, \sigma^* \tilde{\nabla}^2 \Phi) + \dot{\Phi}(\sigma^* \tilde{\nabla}^2 \partial_t \Phi) + \ddot{\Phi}(\partial_t \mathcal{A}, \mathrm{id}) \Phi + \dot{\Phi}(\mathrm{id}) \partial_t \Phi \\ &= \ddot{\Phi}(\partial_t \mathcal{A}, \partial_t \mathcal{A}) + \dot{\Phi}(\sigma^* \tilde{\nabla}^2 \partial_t \Phi) + \dot{\Phi}(\mathrm{id}) \partial_t \Phi \;. \end{split}$$

**Lemma 8.3** (Andrews, [And94, Lemma 3.10]). Let  $f: M^n \times [0,T) \to \mathbb{R}$  and  $\tilde{f}: \mathbb{S}^n \times [0,T) \to \mathbb{R}$  be related by

$$\tilde{f}(\boldsymbol{\nu}(p,t),t) = f(p,t)$$

for all  $p \in M^n$  and  $t \in [0, T)$ . Then

$$\partial_t f = \partial_t \tilde{f} + A^{-1}(\nabla F, \nabla f)$$
.

*Proof.* Differentiating yields

$$\begin{split} \partial_t f &= \partial_t \tilde{f} + \partial_{z_i} \tilde{f} \partial_t \boldsymbol{\nu}^i = \partial_t \tilde{f} + \partial_{p_j} f \partial_{z_i} (\boldsymbol{\nu}^{-1})^j \partial_{p^i} F \\ &= \partial_t \tilde{f} + g_{jk} \partial_{p^k} f a_i^j \partial_{p^i} F = \partial_t \tilde{f} + a_{ij} \partial_{p^i} f \partial_{p^i} F \,, \end{split}$$

where  $(a^{ij})_{ij} = ((h_{ij})_{ij})^{-1}$ .

**Theorem 8.4** (Andrews, [And94, Theorem 5.6]). Suppose X is a strictly convex solution to (25).

(i) If  $\Phi$  is  $\alpha$ -concave for  $0 < \alpha < 1$  ( $\alpha$ -convex for  $\alpha > 1$ ), then

$$\partial_t \Phi + \frac{\alpha \Phi}{(\alpha - 1)t} \le (\ge) 0.$$

for all  $t \in [0,T)$ .

(ii) If  $\Phi$ , is positive and concave (convex), then

$$\sup_{\mathbb{S}^n} (\partial_t \log \Phi) \qquad is decreasing (increasing).$$

*Proof.* We prove the concave cases. For claim (ii), let  $\Phi$  be concave and set  $R := \partial_t \log \Phi$ . Then

$$\partial_t R = \partial_t \left( \frac{\partial_t \Phi}{\Phi} \right) = \frac{\partial_t^2 \Phi}{\Phi} - \frac{(\partial_t \Phi)^2}{\Phi^2}$$

as well as

$$\tilde{\nabla}R = \tilde{\nabla}\left(\frac{\partial_t \Phi}{\Phi}\right) = \frac{\tilde{\nabla}\partial_t \Phi}{\Phi} - \frac{\partial_t \Phi \tilde{\nabla}\Phi}{\Phi^2}$$

and

$$\begin{split} \tilde{\nabla}^2 R &= \tilde{\nabla} \left( \frac{\tilde{\nabla} \partial_t \Phi}{\Phi} - \frac{\partial_t \Phi \tilde{\nabla} \Phi}{\Phi^2} \right) \\ &= \frac{\tilde{\nabla}^2 \partial_t \Phi}{\Phi} - 2 \frac{\tilde{\nabla} \partial_t \Phi \otimes \tilde{\nabla} \Phi}{\Phi^2} - \frac{\partial_t \Phi \tilde{\nabla}^2 \Phi}{\Phi^2} + 2 \frac{\partial_t \Phi (\tilde{\nabla} \Phi)^2}{\Phi^3} \\ &= \frac{\tilde{\nabla}^2 \partial_t \Phi}{\Phi} - 2 \frac{\tilde{\nabla} R \otimes \tilde{\nabla} \Phi}{\Phi} - \frac{\partial_t \Phi \tilde{\nabla}^2 \Phi}{\Phi^2} \,. \end{split}$$

By (28) and (29)

$$\begin{split} \partial_t R &= \frac{1}{\Phi} \left( \dot{\Phi}(\sigma^* \tilde{\nabla}^2 \partial_t \Phi) + \dot{\Phi}(\mathrm{id}) \partial_t \Phi + \ddot{\Phi}(\partial_t \mathcal{A}, \partial_t \mathcal{A}) \right) - \frac{R}{\Phi} \left( \dot{\Phi}(\sigma^* \tilde{\nabla}^2 \Phi) + \dot{\Phi}(\mathrm{id}) \Phi \right) \\ &\leq \dot{\Phi}(\sigma^* \tilde{\nabla}^2 R) + \frac{2}{\Phi} \dot{\Phi} \left( \sigma^* \left( \tilde{\nabla} \Phi \otimes \tilde{\nabla} R \right) \right) + \frac{1}{\Phi^2} \dot{\Phi}(\sigma^* (\partial_t \Phi \tilde{\nabla}^2 \Phi)) \\ &\quad + \frac{1}{\Phi} \dot{\Phi}(\mathrm{id}) \partial_t \Phi - \frac{R}{\Phi} \left( \dot{\Phi}(\sigma^* \tilde{\nabla}^2 \Phi) + \dot{\Phi}(\mathrm{id}) \Phi \right) \\ &= \dot{\Phi}(\sigma^* \tilde{\nabla}^2 R) + \frac{2}{\Phi} \dot{\Phi} \left( \sigma^* \left( \tilde{\nabla} \Phi \otimes \tilde{\nabla} R \right) \right) \; . \end{split}$$

The strong parabolic maximum principle, Theorem C.3, implies (ii), since the first term is an elliptic operator, and the second a gradient term. For claim (i), let  $\Phi$  be  $\alpha$ -concave with  $\alpha < 1$  and set

$$R := t\partial_t \Phi + \frac{\alpha \Phi}{\alpha - 1} \,,$$

which is negative at t = 0. Then

$$\partial_t R = t \partial_t^2 \Phi + \frac{2\alpha - 1}{\alpha - 1} \partial_t \Phi$$

as well as

$$\tilde{\nabla}R = t\tilde{\nabla}\partial_t\Phi + \frac{\alpha}{\alpha - 1}\tilde{\nabla}\Phi$$

and

$$\tilde{\nabla}^2 R = t \tilde{\nabla}^2 \partial_t \Phi + \frac{\alpha}{\alpha - 1} \tilde{\nabla}^2 \Phi \,.$$

By (28), (29) and (27),

$$\begin{split} \partial_t R &= t \left( \dot{\Phi}(\sigma^* \tilde{\nabla}^2 \partial_t \Phi) + \dot{\Phi}(\mathrm{id}) \partial_t \Phi + \ddot{\Phi}(\partial_t \mathcal{A}, \partial_t \mathcal{A}) \right) + \frac{2\alpha - 1}{\alpha - 1} \partial_t \Phi \\ &\leq \dot{\Phi}(\sigma^* \tilde{\nabla}^2 R) - \frac{\alpha}{\alpha - 1} \dot{\Phi}(\sigma^* \tilde{\nabla}^2 \Phi) + t \dot{\Phi}(\mathrm{id}) \partial_t \Phi \\ &+ t \frac{\alpha - 1}{\alpha \Phi} (\dot{\Phi}(\partial_t \mathcal{A}))^2 + \frac{2\alpha - 1}{\alpha - 1} \partial_t \Phi \\ &= \dot{\Phi}(\sigma^* \tilde{\nabla}^2 R) + \frac{\alpha}{\alpha - 1} \left( \dot{\Phi}(\mathrm{id}) \Phi - \partial_t \Phi \right) + t \dot{\Phi}(\mathrm{id}) \partial_t \Phi \\ &+ t \frac{\alpha - 1}{\alpha \Phi} (\partial_t \Phi)^2 + \frac{2\alpha - 1}{\alpha - 1} \partial_t \Phi \\ &= \dot{\Phi}(\sigma^* \tilde{\nabla}^2 R) + \frac{\alpha}{\alpha - 1} \dot{\Phi}(\mathrm{id}) \Phi + t \dot{\Phi}(\mathrm{id}) \partial_t \Phi + t \frac{\alpha - 1}{\alpha \Phi} (\partial_t \Phi)^2 + \partial_t \Phi \\ &= \dot{\Phi}(\sigma^* \tilde{\nabla}^2 R) + \left( \frac{\alpha - 1}{\alpha \Phi} \partial_t \Phi + \dot{\Phi}(\mathrm{id}) \right) \left( t \partial_t \Phi + \frac{\alpha \Phi}{\alpha - 1} \right) \\ &= \dot{\Phi}(\sigma^* \tilde{\nabla}^2 R) + \left( \frac{\alpha - 1}{\alpha \Phi} \partial_t \Phi + \dot{\Phi}(\mathrm{id}) \right) R \,. \end{split}$$

The weak parabolic maximum principle, Theorem C.2, implies that R stays negative as long as the solution exists.

This calculation can easily be transferred to the standard parametrization, by writing the various quantities in terms of the metric and connection on the hypersurface. This is most easily done by considering the change in the evolution equations coming from the modified parametrization. Here we denote by  $A^{-1}$  the map inverse to A.

**Corollary 8.5** (Andrews, [And94, Corollary 5.11]). Suppose X is a strictly convex solution of (25).

(i) If  $\Phi$  is  $\alpha$ -concave for  $\alpha < 1$  ( $\alpha$ -convex for  $\alpha > 1$ ), then

$$\partial_t F - A^{-1}(\nabla F, \nabla F) + \frac{\alpha F}{(\alpha - 1)t} \ge (\le) 0.$$

for all  $t \in [0,T)$ .

(ii) If  $\Phi$  is positive and concave (convex), then  $\sup_{M^n} (\partial_t \log F - FA^{-1}(\nabla \log |F|, \nabla \log |F|)) \quad \text{is decreasing (increasing)}.$ 

**Theorem 8.6** (Andrews, [And94, Theorem 5.17]). Suppose X is a strictly convex solution of (25). The following inequalities apply in the standard parametrization for the cases described, for any points  $p_1, p_2 \in M^n$ , any times  $0 < t_1 < t_2 < T$ , and any curve  $\gamma$  between  $(p_1, t_1)$  and  $(p_2, t_2)$ .

(i) If  $\Phi$  is  $\alpha$ -concave,  $\alpha < 0$ , then

$$\frac{F(p_2,t_2)}{F(p_1,t_1)} \geq \left(\frac{t_1}{t_2}\right)^{\alpha/(\alpha-1)} \exp\biggl(-\frac{1}{4} \int_{\gamma} F^{-1} A(\dot{\gamma},\dot{\gamma}) \, dt \biggr) \; .$$

(ii) If  $\Phi$  is  $\alpha$ -convex,  $\alpha > 1$ , then

$$\frac{F(p_2, t_2)}{F(p_1, t_1)} \ge \left(\frac{t_1}{t_2}\right)^{\alpha/(\alpha - 1)} \exp\left(-\frac{1}{4} \int_{\gamma} |F|^{-1} A(\dot{\gamma}, \dot{\gamma}) dt\right).$$

(iii) If  $\Phi$  is convex and positive, then

$$\frac{F(p_2, t_2)}{F(p_1, t_1)} \ge \exp(-C(t_2 - t_1)) \exp\left(-\frac{1}{4} \int_{\gamma} |F|^{-1} A(\dot{\gamma}, \dot{\gamma}) dt\right) ,$$

where  $C = \lim_{t \searrow 0} \sup_{M^n} \left( \partial_t \log |F| - FA^{-1}(\nabla \log |F|, \nabla \log |F|) \right),$ 

*Proof.* Along a curve  $\gamma$ ,

$$D_{\dot{\gamma}} \log F = \partial_t \log F + \langle \dot{\gamma}, \nabla \log F \rangle.$$

By Corollary 8.5(i), Cauchy-Schwarz and Young,

$$D_{\dot{\gamma}} \log F \ge F A^{-1} (\nabla \log F, \nabla \log F) + \langle \dot{\gamma}, \nabla \log F \rangle - \frac{\alpha}{(\alpha - 1)t}$$
$$\ge -\frac{1}{4} F^{-1} A(\dot{\gamma}, \dot{\gamma}) - \frac{\alpha}{(\alpha - 1)t}.$$

Integrating along  $\gamma$  yields claim (i). For claim (ii),

$$D_{\dot{\gamma}} \log F \ge C - \frac{1}{4} F^{-1} A(\dot{\gamma}, \dot{\gamma}) \ge -C - \frac{1}{4} F^{-1} A(\dot{\gamma}, \dot{\gamma}).$$

**Theorem 8.7** (Hamilton [Ham95b, Theorem 1.3]). Let  $X: M^n \times (-\infty, T) \to \mathbb{R}^{n+1}$  be an ancient mean curvature flow of a complete, strictly convex hypersurface with bounded second fundamental form at every time and such that H takes its maximum in space and time. Then, X is a translating flow.

*Proof.* Define

$$Z := \partial_t H + \frac{H}{2(t - t_0)} - A^{-1}(\nabla H, \nabla H)$$

then

$$(\partial_t - \Delta)Z = 2g^{ij}a^{kl}J_{ik}J_{jl} + \left(|A|^2 - \frac{2}{t - t_0}\right)Z \ge \left(|A|^2 - \frac{2}{t - t_0}\right)Z$$

where

$$J_{ik} = \nabla_{ik}^2 H + H h_{ik}^2 - a^{sr} \nabla_s H \nabla_r h_{ik} + \frac{h_{ik}}{2(t - t_0)}$$

By the weak maximum principle,  $Z \ge 0$ . On an eternal solution which attains its maximum in space and time, Z = 0 at this maximum when sending  $t_0 \to -\infty$ . Hence  $Z \equiv 0$  and

$$\partial_t H = A^{-1}(\nabla H, \nabla H)$$
.

Since

$$g^{ik} = g^{kl}\delta^i_l = g^{kl}h_{jl}a^{ij} = h^k_i a^{ij}$$

and, by Codazzi and  $a^{ij}h_{jk} = \delta_k^i$ ,

$$\begin{split} a^{il}\nabla_l H &= a^{il}g^{km}\nabla_l h_{km} = a^{il}g^{km}\nabla_k h_{lm} \\ &= -a^{il}g^{km}h_{ls}h_{mj}\nabla_k a^{sj} = -h_i^k\nabla_k a^{ij} \,, \end{split}$$

we obtain

$$\begin{split} 0 &= -a^{il}\nabla_l H \nabla_i H + \Delta H + H |A|^2 \\ &= \left(\nabla_k a^{ij}\nabla_i H + a^{ij}\nabla_k \nabla_i H + H h_k^j\right) h_j^k \,. \end{split}$$

Consider the vector

$$V = a^{ij} \nabla_i H \nabla_j X + H \boldsymbol{\nu} \,.$$

Since

$$\nabla_k \nabla_j X = \langle \partial_k \partial_j X, \boldsymbol{\nu} \rangle \boldsymbol{\nu} = -h_{jk} \boldsymbol{\nu}$$

and

$$\nabla_k \boldsymbol{\nu} = h_k^j \nabla_j X = g^{ij} h_{ik} \nabla_j X$$

as well as  $a^{ij}h_{jk} = \delta^i_k$ , we obtain

$$\nabla_k V = \nabla_k a^{ij} \nabla_i H \nabla_j X + a^{ij} \nabla_k \nabla_i H \nabla_j X + a^{ij} \nabla_i H \nabla_k \nabla_j X + \nabla_k H \boldsymbol{\nu} + H \nabla_k \boldsymbol{\nu}$$
$$= \left( \nabla_k a^{ij} \nabla_i H + a^{ij} \nabla_k \nabla_i H + H h_k^j \right) \nabla_j X + \left( \nabla_k H - a^{ij} \nabla_i H h_{jk} \right) \boldsymbol{\nu} = 0.$$

On the other hand, at a fixed point so that the Christoffel symbols vanish,

$$\partial_t a^{ij} = -a^{ik} a^{jl} \partial_t h_{kl} = -a^{ik} a^{jl} \left( \nabla_k \nabla_l H - H g^{ms} h_{lm} h_{ks} \right)$$
$$= -a^{ik} a^{jl} \nabla_k \nabla_l H + H g^{ij}.$$

and

$$\partial_t \partial_i H = \partial_i \left( a^{kl} \partial_k H \partial_l H \right) = -H h_i^l \partial_l H + a^{kl} \partial_k H \partial_i \partial_l H$$

as well as

$$\partial_t \partial_j X = -\partial_j (H \boldsymbol{\nu}) = -\partial_j H \boldsymbol{\nu} - H h_j^k \partial_k X$$

and

$$\partial_t \mathbf{\nu} = q^{ij} \partial_i H \partial_i X$$

Together, we obtain,

$$\begin{split} \partial_t V &= \partial_t a^{ij} \partial_i H \partial_j X + a^{ij} \partial_t \partial_i H \partial_j X + a^{ij} \partial_i H \partial_t \partial_j X + \partial_t H \boldsymbol{\nu} + H \partial_t \boldsymbol{\nu} \\ &= \left( H g^{ij} \partial_i H - a^{ik} a^{jl} \nabla_k \nabla_l H \partial_i H - a^{ij} H h^l_i \partial_l H + a^{ij} a^{kl} \partial_k H \partial_i \partial_l H \right. \\ &\quad \left. + H g^{ij} \partial_i H \right) \partial_j X - a^{ij} \partial_i H H h^k_i \partial_k X + \left( a^{kl} \partial_k H \partial_l H - a^{ij} \partial_i H \partial_j H \right) \boldsymbol{\nu} = 0 \,. \end{split}$$

Hence V is a constant vector field in space and time. Let  $t_1 \in (-\infty, T)$  and  $\phi: M^n \to M^n$  be a diffeomorphism with  $\phi(\cdot, t_1) = \mathrm{id}$  and

$$\partial_t \phi = -a^{ij} \nabla_i H \nabla_j X$$

and  $\tilde{X}(p,t) = X(\phi(p,t),t)$ . By Theorem 1.3,  $\tilde{X}(M^n,t) = X(\phi(M^n,t),t) = M_t$  and

$$\tilde{X}(p,t) - \tilde{X}(p,t_1) = X(\phi(p,t),t) - X(p,t_1) = \int_{t_1}^{t} \langle DX, \partial_t \phi \rangle + \partial_t X \, d\tau$$
$$= -\int_{t_1}^{t} a^{ij} \nabla_i H \nabla_j X + H \nu \, d\tau = -(t-t_1) V$$

so the  $M_t = M_{t_1} - (t - t_1)V$  and the surfaces move by translation in direction of -V.

### 9. Noncollapsing

#### 10. Convexity estimates

# 11. Type-II singularities (continued)

#### 12. Two-convex hypersurfaces

## Appendix A. Hypersurfaces in $\mathbb{R}^{n+1}$

A topological space is called Hausdorff space if for any two distinct points there exists a neighbourhood of each which is disjoint from the neighbourhood of the other. A topological space  $M^n$  is called locally Euclidean of dimension n, if  $M^n$  can be covered with open sets where every set is homeomorphic to an open subset of  $\mathbb{R}^n$ . A pair  $(U,\varphi)$ , where  $U \subset M^n$  is open and  $\varphi: U \to \varphi(U) \subset \mathbb{R}^n$  is a homeomorphism, is called chard of  $M^n$ . A collection A of chards is called atlas of  $M^n$  if

$$M^n \subset \bigcup_{(U,\varphi)\in A} U$$
.

Two chards  $(U, \varphi)$  and  $(V, \psi)$  are called  $C^k$ -compatible,  $k \geq 1$ , if

$$\psi \circ \varphi^{-1} : \varphi(U \cap V) \to \psi(U \cap V)$$

is a  $C^k$ -diffeomorphism. An atlas is called of class  $C^k$ , if each of its chards are  $C^k$ -compatible. If A is a  $C^k$ -atlas, there exists exactly one maximal  $C^k$ -atlas  $A_0$  with  $A \subset A_0$ ; it contains all chards which are  $C^k$  compatible with the chards of A. A differentiable  $(C^k$ -)structure on  $M^n$  is a maximal  $C^k$ -atlas on  $M^n$ . A local Euclidean Hausdorff space with a differentiable structure is called differentiable manifold.

Let  $N^{n+m}$  be a differentiable manifold. A subset  $M^n \subset N^{n+m}$ ,  $n,m \geq 1$ , is called n-dimensional  $C^k$ -submanifold of  $N^{n+m}$  if for every  $x \in M^n$  there exists an open neighbourhood  $U \subset N^{n+m}$  and a  $C^k$  diffeomorphism  $\varphi: U \to \varphi(U) \subset \mathbb{R}^{n+m}$  with

$$\varphi(U \cap M) = \varphi(U) \cap (\mathbb{R}^n \times \{0_{\mathbb{R}^m}\}).$$

Such an  $M^n$  owns a  $C^k$ -atlas, that is

$$A := \{(U \cap M, \varphi|_{U \cap M}) \mid \text{ where } (U, \varphi) \text{ as above} \}.$$

Then,  $M^n$  is locally Euclidean of dimension m and

$$(\psi|_{V\cap M}) \circ (\varphi|_{U\cap M})^{-1} = \psi \circ \varphi^{-1}|_{(\mathbb{R}^n \times \{0\}) \cap \varphi(U\cap V)} \in C^k$$

for two diffeomorphisms  $\psi$  and  $\varphi$ .

A topological manifold with boundary is a Hausdorff space in which every point has a neighborhood homeomorphic to an open subset of the Euclidean half-space  $\mathbb{R}^n_+ = \{(x_1, \dots, x_n) \in \mathbb{R}^n \mid x_n \geq 0\}$ . The boundary  $\partial M^n$  of  $M^n$  is the set of all points  $p \in M^n$  such that  $(\varphi(p))^n = 0$  for all chards  $(U, \varphi)$  of  $M^n$ . If  $M^n$  is a manifold with boundary, then the interior int  $M^n = M^n \setminus \partial M^n$  is a manifold (without boundary) of dimension n and boundary  $\partial M^n$  is a manifold (without boundary) of dimension n-1.

Let  $M^n$  be an abstract, smooth, compact, n-dimensional manifold without boundary and X a smooth immersion (rank  $DX \equiv n$ ) with

$$X: M^n \to \mathbb{R}^{n+m}$$
.

We call  $M := X(M^n)$  a hypersurface in  $\mathbb{R}^{n+m}$ . For all  $p \in M^n$  and  $v, w \in T_pM^n$ , the embedding X induces an isomorphism

$$dX_p: T_pM^n \to T_{X(p)}M$$
,

and the first fundamental form or metric  $g_p: T_pM^n \times T_pM^n \to \mathbb{R}$  with

$$g_p(v,w) := \langle dX_p(v), dX_p(w) \rangle_{\mathbb{R}^{n+m}}$$
.

Let  $(U_i, \varphi_i)_{i \in I}$  be an atlas of  $M^n$  and

$$\partial_i = \frac{\partial}{\partial p_i} = d\varphi^{-1}(e_i) \in TM^n$$

then the matrix entries of the metric are

$$g_{ij} = g(\partial_i, \partial_j) = \langle dX(\partial_i), dX(\partial_j) \rangle_{\mathbb{R}^{n+m}} = \langle \partial_i X, \partial_j X \rangle_{\mathbb{R}^{n+m}} = \delta_{\alpha\beta} \partial_i X^{\alpha} \partial_j X^{\beta}$$

for  $1 \leq \alpha, \beta \leq n + m$ . We define by  $(g^{ij})_{ij}$  the coordinate dependent inverse of the matrix  $(g_{ij})_{ij}$  and the measure

$$d\mu^n = \sqrt{\det(g_{ij})} \, dp \, .$$

Observe that

$$\partial_k g_{ij} = \langle \partial_k \partial_i X, \partial_j X \rangle + \langle \partial_i X, \partial_k \partial_j X \rangle$$

and

$$\partial_k q^{ij} = -q^{pi} q^{qj} \partial_k q_{pq}$$
.

The corresponding Levi-Cevita connection on  $M^n$  is given by

$$\nabla_v w = dX^{-1} \Big( \Big( D_{dX(v)} dX(w) \Big)^{\top} \Big) .$$

Here D is the standard connection in  $\mathbb{R}^{n+m}$ , and  $^{\top}$  denotes the tangential component with respect to M, that is the orthogonal projection onto  $dX(p)(T_pM^n) = T_{X(p)}M$ . The connection can be evaluated in coordinates in terms of the *Christoffel symbols*  $\Gamma^k_{ij}$  defined by

$$\nabla_{\partial_i}\partial_i = \Gamma^k_{ij}\partial_k \,,$$

where  $\Gamma^k_{ij}$  is explicitly given by We define the Christoffel symbols by

$$\Gamma_{ij}^k := g^{kl} \langle \partial_i \partial_j X, \partial_l X \rangle.$$

Here and in the following, we sum over repeated indices. Then,

$$\Gamma_{ij}^k \partial_k X = \langle \partial_i \partial_j X, \partial_l X \rangle \partial_l X.$$

At a fixed point, we can choose a coordinate system such that  $\Gamma_{ij}^k = 0$ . We calculate

$$0 = \partial_k \delta^i_j = \partial_k (g^{il} g_{jl}) = g^{il} \partial_k g_{jl} + g_{jl} \partial_k g^{il},$$

so that

$$\begin{split} \partial_k g^{ij} &= -g^{il} g^{jm} \partial_k g_{lm} = -g^{il} g^{jm} \partial_k \langle \partial_l X, \partial_m X \rangle \\ &= -g^{il} g^{jm} (\langle \partial_k \partial_l X, \partial_m X \rangle + \langle \partial_l X, \partial_k \partial_m X \rangle) = -g^{il} \Gamma^j_{kl} - g^{jm} \Gamma^i_{km} \,. \end{split}$$

Being in a Levi–Cevita connection the *Lie bracket*  $[\cdot, \cdot]$  is given by

$$[v, w] = \nabla_v w - \nabla_w v = (v(\mu^k) - w(\lambda^k)) \partial_k.$$

The tangential gradient of a function  $f \in C^1(M)$  is given by

$$\nabla^M f = g^{ij} \partial_i f \partial_j .$$

The tangential divergence  $\operatorname{div}_M: T_pM^n \to \mathbb{R}$  is given by

$$\operatorname{div}_{M} v = g^{ij} \langle \partial_{i} v, \partial_{j} X \rangle_{\mathbb{R}^{n+m}}.$$

For the embedding vector X, we therefore have

$$\operatorname{div}_M X = g^{ij} \langle \partial_i X, \partial_j X \rangle_{\mathbb{R}^{n+m}} = g^{ij} g_{ij} = n.$$

For  $\omega = df = \frac{\partial f}{\partial p_i} dp^i$ , we obtain the Hessian of the function f

$$(\operatorname{Hess}_M f)(v, w) := (\nabla^2 f)(v, w)$$

or in coordinates

$$\nabla_i \nabla_j f = (\operatorname{Hess}_M f)(\partial_i, \partial_j) = \partial_i \partial_j f - \Gamma_{ij}^k \partial_k f.$$

The Laplace–Beltrami operator  $\Delta_M: C^2(M^n) \to C^0(M^n)$  is defined as

$$\Delta_M f := \frac{1}{\sqrt{\det g_{kl}}} \partial_j \Big( \sqrt{\det g_{kl}} \, g^{ij} \partial_j f \Big) = \operatorname{div}_M(\nabla^M f) = g^{ij} \nabla_i \nabla_j f \,.$$

We define the second fundamental form  $\mathbf{A}_p: T_pM^n \times T_pM^n \to \left(T_{X(p)}M\right)^{\perp}$  by

$$\mathbf{A}_{p}(v,w) := -\sum_{k=1}^{m} \left\langle D_{dX_{p}(v)} dX_{p}(w), \boldsymbol{\nu}_{k}(p) \right\rangle \boldsymbol{\nu}_{k}(p)$$

$$= \sum_{k=1}^{m} \left\langle dX_{p}(w), D_{dX_{p}(v)} \boldsymbol{\nu}_{k}(p) \right\rangle \boldsymbol{\nu}_{k}(p),$$

where  $\{\nu_k\}_{1\leq k\leq m}$  is an orthonormal frame for  $(TM)^{\perp}$ . In coordinates  $\{p_i\}_{1\leq i\leq n}$ ,

$$\mathbf{A}_{ij} := \mathbf{A}_p(\partial_i, \partial_j) = \sum_{k=1}^m \langle \partial_i X, \partial_j \boldsymbol{\nu}_k \rangle \, \boldsymbol{\nu}_k \,.$$

The mean curvature vector  $\mathbf{H}: M \to (TM)^{\perp}$  is the trace of the second fundamental form

$$\mathbf{H} := -g^{ij}\mathbf{A}_{ij} = -g^{ij}\sum_{k=1}^{m} \langle \partial_i X, \partial_j \boldsymbol{\nu}_k \rangle \, \boldsymbol{\nu}_k = -\sum_{k=1}^{m} \operatorname{div}(\boldsymbol{\nu}_k) \boldsymbol{\nu}_k \,.$$

We calculate that

$$\Delta_{M}X = g^{ij} \left( \partial_{i} \partial_{j} X - \Gamma_{ij}^{k} \partial_{k} X \right) = g^{ij} \sum_{k=1}^{m} \left\langle \partial_{i} \partial_{j} X, \boldsymbol{\nu}_{k} \right\rangle \boldsymbol{\nu}_{k}$$
$$= -g^{ij} \sum_{k=1}^{m} \left\langle \partial_{i} X, \partial_{j} \boldsymbol{\nu}_{k} \right\rangle \boldsymbol{\nu}_{k} = \mathbf{H}.$$

For a submanifold  $\Sigma$  of M, the mean curvature vector is given by

$$\mathbf{H}_{\Sigma} = -\sum_{k=1}^{m} \operatorname{div}_{\Sigma}(\boldsymbol{\nu}_{k}) \boldsymbol{\nu}_{k} - \operatorname{div}_{\Sigma}(\boldsymbol{\nu}_{\Sigma}) \boldsymbol{\nu}_{\Sigma},$$

where  $\nu_{\Sigma}$  is the unit co-normal of  $\Sigma$ . Since  $\nu_{\Sigma}$  tangential to M,

$$\langle \mathbf{H}_{\Sigma}, \boldsymbol{\nu}_{\Sigma} \rangle = -\operatorname{div}_{\Sigma} \boldsymbol{\nu}_{\Sigma}$$

and on  $\Sigma$ ,

$$\begin{split} \Delta_{\Sigma} X &= g_{\Sigma}^{ij} \left( \partial_{i} \partial_{j} X - \Gamma_{ij}^{k} \partial_{k} X \right) \\ &= \sum_{k=1}^{m} g_{\Sigma}^{ij} \left\langle \partial_{i} \partial_{j} X, \boldsymbol{\nu}_{k} \right\rangle \boldsymbol{\nu}_{k} + g_{\Sigma}^{ij} \left\langle \partial_{i} \partial_{j} X, \boldsymbol{\nu}_{\Sigma} \right\rangle \boldsymbol{\nu}_{\Sigma} = \mathbf{H}_{\Sigma} \,. \end{split}$$

For m = 1,

$$\mathbf{A}(v,w) = A(v,w)\boldsymbol{\nu}\,,$$

where  $\nu$  is the outward pointing unit normal to M and  $A:TM^n\times TM^n\to\mathbb{R}$  is given by

$$A(v,w) = -\langle D_{dX(v)}dX(w), \boldsymbol{\nu} \rangle = \langle dX(w), D_{dX(v)}\boldsymbol{\nu} \rangle.$$

where  $\nu$  is the outward pointing unit normal to M. In coordinates,

$$h_{ij} := A(\partial_i, \partial_j) = -\langle \partial_i \partial_j X, \boldsymbol{\nu} \rangle = \langle \partial_i X, \partial_j \boldsymbol{\nu} \rangle.$$

Let  $\lambda_1 \leq \lambda_2 \leq \cdots \leq \lambda_n$  be the eigenvalues of A, that is

$$h_{ij}\xi_k^i\xi_k^j = \lambda_k g_{ij}$$

for eigenvectors  $\xi_k \in TM$  and k = 1, ..., n. The Weingarten operator  $S: TM^n \to TM^n$  is given by

$$S(v) := dX^{-1}(D_{dX(v)}\boldsymbol{\nu})$$

so that

$$A(v, w) = g(v, S(w)),$$

where in coordinates,

$$h_j^i := g^{ik} h_{kj}$$

and the Weingarten equations by

$$\partial_i \boldsymbol{\nu} = h_i^j \partial_j X \,. \tag{30}$$

The norm of the second fundamental form is given by

$$|A|^2 = g^{ik}g^{lj}h_{kl}h_{ij} = h^{ij}h_{ij} \,,$$

and the mean curvature vector is given by

$$\mathbf{H} = -g^{ij}h_{ij}\boldsymbol{\nu} = -H\boldsymbol{\nu}\,,$$

where we define the  $mean\ curvature\ H$  of M as the trace of the second fundamental form with

$$H = g^{ij} h_{ij} = \operatorname{div}_M \boldsymbol{\nu}$$
.

We have the Gauss formula

$$\nabla_i \nabla_j X = \partial_i \partial_j X - \Gamma_{ij}^k \partial_k X = -h_{ij} \nu \tag{31}$$

which as before leads to  $\Delta_M X = \mathbf{H}$ . More useful identities are the Codazzi equations in  $\mathbb{R}^{n+1}$ 

$$\nabla_k h_{ij} - \nabla_j h_{ik} = \Gamma^l_{ij} h_{lk} - \Gamma^l_{ik} h_{lj} \tag{32}$$

and Simons' identity

$$\Delta h_{ij} = \nabla_i \nabla_j H + H h_{ik} h_i^k - |A|^2 h_{ij}. \tag{33}$$

We define the Riemannian curvature tensor as

$$R_{lij}^k := \nabla_i \Gamma_{jl}^k - \nabla_j \Gamma_{il}^k + \Gamma_{im}^k \Gamma_{jl}^m - \Gamma_{jm}^k \Gamma_{il}^m.$$

Moreover, we set

$$R_{klij} := g^{kr} R_{lij}^r$$

and define the Ricci tensor by

$$R_{ik} := R_{ijkl}g^{jl}$$

and the scalar curvature by

$$R := R_{ij}g^{ij}$$
.

The Gauss equation

$$R_{ijkl} = h_{ik}h_{jl} - h_{il}h_{jk} \tag{34}$$

**Theorem A.1** (First variation of the area formula, see [Sim83, p. 51]). Let  $M \subset \mathbb{R}^{n+1}$  be a smooth, compact, n-dimensional hypersurface with boundary. Let  $U \subset \mathbb{R}^{n+1}$  be a open and bounded such that  $M \subset U$ . Let  $\phi: U \times (-1,1) \to U$  be a one-parameter family of  $C^{2;1}$ -diffeomorphisms. Set  $M_t := \phi(M,t)$  and  $v(p) := \partial_t \phi(p,0)$ . Then

$$\partial_t|_{t=0}\mu^n(M_t) = \int_M \operatorname{div}_M v \, d\mu^n.$$

**Theorem A.2** (Divergence theorem, see [Sim83, p. 43], [DHTK10, p. 304], [Eck04, p. 116]). Let  $M \subset \mathbb{R}^{n+1}$  be a smooth, compact, n-dimensional manifold with boundary. Let v be a  $C^1$ -vector field on M. Then

$$\int_{M} \operatorname{div}_{M} v \, d\mu^{n} = -\int_{M} \langle v, \mathbf{H}_{M} \rangle_{\mathbb{R}^{n+1}} \, d\mu^{n} + \int_{\partial M} \langle v, \boldsymbol{\nu}_{\partial M} \rangle_{\mathbb{R}^{n+1}} \, d\mu^{n-1} \, .$$

Proof.

**Theorem A.3** (Rademachers theorem, see [Fed69, Theorem 3.1.6]). Let  $U \subset \mathbb{R}^n$  be open and  $f: U \to \mathbb{R}^m$  be Lipschitz continuous. Then f is differentiable almost everywhere in U.

# APPENDIX B. SARD'S THEOREM

Section copied from [Sch05, Section 3]. See also [BJ73].

**Definition B.1.** Let  $f: M \to N$  differentiable. A point  $p \in M$  is called *regular*, if the differential of f in p is surjektiv. A point  $q \in N$  is called *regular value*, if  $f^{-1}(q)$  consists of regular points. Non-regular points or values are called *singular* or *critical*.

We want to prove the following theorem.

**Theorem B.2** (Sard's theorem). Let  $M^m$  and  $N^n$  be differentiable manifolds with a countable basis of their topology. The critical set S of a  $C^k$  function  $f: M \to N$  consists of those points at which the differential  $df: TM \to TN$  has rank less than n as a linear transformation. If  $k \ge \max\{n-m+1,1\}$ , then the image of S has Lebesgue measure zero as a subset of N.

**Corollary B.3.** Let  $M^m$  be a differentiable manifold and  $f: M^m \to \mathbb{R}^n$  a differentiable. Then  $f^{-1}(x) \subset M^m$  is a differentiable submanifold of co-dimension n for almost every  $x \in \mathbb{R}^n$ .

**Remark B.4.** The set  $f^{-1}(x)$  can be empty. Sard's theorem also holds for maps  $f: \mathbb{R}^n \to \mathbb{R}^p$ ,  $f \in C^k$  with  $k > \max\{n-p,0\}$  and manifolds with according dimensions.

**Definition B.5.** A subset  $C \subset \mathbb{R}^n$  is of *measure zero*, if for every  $\varepsilon > 0$  there exists a sequence  $(W_i)_{i \in \mathbb{N}}$  of cubes in  $\mathbb{R}^n$  with

$$C \subset \bigcup_{i \in \mathbb{N}} W_i$$
 and  $\sum_{i \in \mathbb{N}} |W_i| < \varepsilon$ .

Remark B.6. (i) The countable set of zero sets is again a zero set.

(ii) One obtains an equivalent definition for open oder closed cubes or balls.

**Lemma B.7.** Let  $U \subset \mathbb{R}^m$  be open and  $C \subset U$  of measure zero. Let  $f: U \to \mathbb{R}^m$  be Lipschitz. Then f(C) has measure zero.

Proof. Exercise. 
$$\Box$$

**Definition B.8.** A subset C of a differentiable manifold has measure zero, if for every chard  $h: U \to U' \subset \mathbb{R}^m$  the set  $h(C \cap U) \subset \mathbb{R}^m$  is of measure zero.

**Remark B.9.** The assumption of differentiability is important here, since zero sets are not necessarily maintained under homeomorphisms. Since a manifold owns a countable basis of the topologie, there exists an atlas with countably many chards. It is sufficient to apply the definition for such chards. Well-definedness follows, since zero sets are maintained under differentiable chard changes and countable unions.

**Lemma B.10.** An open covering of the interval [0,1] by subintervals contains a countable cover  $[0,1] = \bigcup_{j=1}^k I_j$  with  $\sum_{j=1}^k |I_j| \leq 2$ .

*Proof.* Due to the compactness, there exists a finite subcover. Choose one where no interval can be left out without loosing the covering property. Let the intervals  $I_j$ , j = 1, ..., k be numbered so that with  $I_j = (a_j, b_j)$  always holds  $a_j < a_{j+1}$ , j = 1, ..., k-1. Minimality and covering property imply  $a_i < a_{i+1} < b_i < a_{i+2}$ . So that

$$\sum_{i} (b_{i} - a_{i}) = \sum_{i} (a_{i+1} - a_{i}) + \sum_{i} (b_{i} - a_{i+1})$$

$$< \sum_{i} (a_{i+1} - a_{i}) + \sum_{i} (a_{i+1} - a_{i+1}) \le 2,$$

where we used that we have telecope sums in the end.

**Theorem B.11** (Fubini). Let  $\mathbb{R}^{n-1}_t := \{x \in \mathbb{R}^n \mid x^n = t\} \subset \mathbb{R}^n$ . Let  $C \subset \mathbb{R}^n$  be compact and  $C_: = C \cap \mathbb{R}^{n-1}_t$  be of measure zero in  $\mathbb{R}^{n-1}_t \cong \mathbb{R}^{n-1}$  for all  $t \in \mathbb{R}$ . Then C is of measure zero in  $\mathbb{R}^n$ .

*Proof.* Since the property of being of measure zero is maintained under countable unions, we can assume that  $C \subset \mathbb{R}^{n-1} \times [0,1]$ . For  $t \in [0,1]$ ,  $C_t$  is of measure zero in  $\mathbb{R}^{n-1} \times \{t\}$ . Let  $\varepsilon > 0$  and  $W_t^i$  be a cover of  $C_t$  by open cubes with  $\sum_i |W_t^i| < \varepsilon$ . Define  $W_t := \bigcup_i W_t^i$  identify these with subsets of  $\mathbb{R}^{n-1}$ . The function  $|x^n - t|$  is for fixed  $t \in [0,1]$  on C continuous, vanishes exactly on  $C_t$  und attains a positive minimum in the compact set  $C \setminus (W_t \times [0,1])$ , which we call  $\alpha$ . It follows

$$\{x \in C : |x^n - t| < \alpha\} \subset W_t \times I_t^{\alpha},$$

where  $I_t^{\alpha} = (t - \alpha, t + \alpha)$  and  $\bigcup_t I_t^{\alpha} = [0, 1]$ . Choose a subcover of [0, 1] among the intervals  $I_t^{\alpha}$  with  $\sum_{t_i} |I_{t_i}^{\alpha}| \leq 2$ . Observe that  $\alpha = \alpha(t_i)$ . It holds

$$C \subset \bigcup_{t_j,i} W_{t_j}^i \times I_{t_j}^\alpha \,,$$

where i is the index of the cube and we take the union over cuboids. Moreover,

$$\sum_{t_j,i} |W_{t_j}^i \times I_{t_j}^{\alpha}| \le 2\varepsilon.$$

Sending  $\varepsilon \to 0$  yields the lemma.

**Remark B.12.** The requirement that C is compact, can be weakened as follows: C is a countable union of compact sets, that each suffice the assumptions of the theorem. This is fulfilled by closed and open sets (which cannot be zero sets), for images of these set under continuous maps, countable union und finite intersections of these.

Proof of Theorem B.2. After introducing maps it is sufficient to show: Let  $U \subset \mathbb{R}^n$  be open,  $f: U \to \mathbb{R}^p$  smooth and  $D \subset U$  be the set of critical points of f, then  $f(D) \subset \mathbb{R}^p$  has measure zero.

We prove by induction over n. In case n = 0,  $\mathbb{R}^n$  is a point. So, f(U) is at most a point and has measure zero. Assume the claim is true for the case n-1. We proof the case n. Let  $D_i \subset U$  be the set of all points points, in which the partial derivative of order  $\leq i$  vanish. We obtain the decreasing sequence of relatively closed sets

$$D\supset D_1\supset D_2\supset\ldots$$

We claim that

- (i)  $f(D \setminus D_1)$  is of measure zero,
- (ii)  $f(D_i \setminus D_{i+1})$  is of measure zero,
- (iii) for k big enough,  $f(D_k)$  is of measure zero.

We observe, that (iii) is neccessary, since also the points, in which all derivatives vanish, can be captured. By Remark B.12, all sets occuring in (i)–(iii) can be used. Moreover, it is sufficient to prove that every point in  $D \setminus D_1$  resp.  $D_i \setminus D_{i+1}$  resp.  $D_k$  has a neighbourhood V, so that  $f(V \cap (D \setminus D_1))$  resp.  $f(V \cap (D_i \setminus D_{i+1}))$  resp.  $f(V \cap D_k)$  are of measure zero. The claim then follows, since the countable union of zero set is again a zero set.

Proof of (i): Assume, that  $p \geq 2$ , since for p = 1 we already have  $D = D_1$ . Let  $x_0 \in D \setminus D_1$ . Since  $x_0 \notin D_1$ , there exists a partial derivative that is not vanishing in  $x_0$ , w.l.o.g.  $\partial_1 f \neq 0$ . Define  $h: U \to \mathbb{R}^n$  by

$$h: x = (x^1, \dots, x^n) \mapsto (f^1(x), x^2, \dots, x^n)$$

Then h is not singular in  $x_0$ . Hence there exists a neighbourhood V of  $x_0$ , so that  $h: V \to h(V) = V'$  is a diffeomorphism. Define  $g:= f \circ h^{-1}$ . In a neighbourhood of h(x), g is of the form

$$g:(z^1,\ldots,z^n)\mapsto (z^1,g^2(z),\ldots,g^n(z)).$$

The hyperplane  $\{z \,|\, z^1=t\}$  is (locally) mapped into the hyperplane  $\{y \,|\, y^1=t\}$ . Define

$$g_t: \{t\} \times \mathbb{R}^{n-1} \cap V' \to \{t\} \times \mathbb{R}^{p-1}$$

als restriction of g. We have

$$Dg_t = \begin{pmatrix} 1 & 0 \\ ? & Dg \end{pmatrix} .$$

Hence a point in  $(\{t\} \times \mathbb{R}^{n-1}) \cap V'$  is critical for g if and only if it is for  $g_t$ . By the induction assumption the set of critical values of  $g_t$  is of measure zero in  $\{t\} \times \mathbb{R}^{p-1}$ . Since g maps entsprechende hyperplanes onto itself, the set of critical values of g also has a intersection of measure zero with the hyperplane  $\{y \mid y^1 = t\}$ . By Fubini, Theorem B.11, the critical values of g have measure zero. Since f and g only differ by an diffeomorphism, also the criticalen values of f have measure zero. This holds locally, as long as  $\partial_1 f \neq 0$ . This proves (i).

Proof of (ii): We argument similarly as in the proof of (i). Let  $x_0 \in D_k \setminus D_{k+1}$ . Then there exist a non-vanishing (k+1)-st derivative, w.l.o.g.

$$\frac{\partial^{k+1} f^1}{\partial x^1 \partial x^{\nu_1} \dots \partial x^{\nu_k}}(x_0) \notin 0.$$

Assume, that this holds in a neighbourhood V of  $x_0$ . Define  $w: V \to \mathbb{R}$  by

$$w:=\frac{\partial^k f^1}{\partial x^{\nu_1}\dots\partial x^{\nu_k}}(x_0)\neq 0.$$

It holds w(x) = 0,  $\frac{\partial}{\partial x^1} w(x) \neq 0$ . The map

$$h: x \to (w(x), x^2, \dots, x^n)$$

defines a diffeomorphism  $h: V \to V' = h(V)$ . w and therefore all k-th derivatives of  $f^1$  vanish at most for  $x = x_0$ . Hence

$$h(D_k \cap V) \subset \{0\} \times \mathbb{R}^{n-1} \subset \mathbb{R}^n$$
.

Define

$$q: f \circ h^{-1}: V' \to \mathbb{R}^p$$

and

$$g_0: \{0\} \times \mathbb{R}^{n-1} \cap V' \to \mathbb{R}^p$$
.

By the induction assumption, the set of critical values of  $g_0$  has measure zero. Let  $x \in h(D_k \cap V)$ . Then all derivatives of g up to order k vanish there. Since  $h(D_k \cap V) \subset \{0\} \times \mathbb{R}^{k-1}$ ,  $g_0$  is defined there and has vanishing derivatives up to order k. In particular, all first derivatives vanish there as well and thus we are dealing with critical points of  $g_0$ . Hence

$$(g_0 \circ h)(D_k \cap V) = (g \circ h)(D_k \cap V) = f(D_k \cap V)$$

has measure zero.

Proof of (iii): The set U ist countable union of cubes. Let  $W \subset U$  be a cube with side length  $a \leq 1$  and let k > n - 1. It is sufficient to show, that  $f(W \cap D_k)$  is of measure zero. By Taylor it holds that

$$f(x+h) = f(x) + R(x,h)$$

with

$$|R(x,h)| \le c|h|^{k+1}$$

for  $x \in D_k \cap W$  and  $x + h \in W$ , where the constant c only depends on f and W. We devide W in  $r^n$  cubes with side length a/r,  $r \in \mathbb{N}$ . If  $W_1$  is a cube of this partitioning, which contains a point  $x \in D_k$ , then every other point in  $W_1$  can be described as x + h with  $|h| \leq \sqrt{n}a/r$ . Hence with Taylor

$$|f(x+h) - f(x)| \le c \left(\frac{\sqrt{na}}{r}\right)^{k+1}$$
.

So that  $f(W_1)$  is contained in a cube with side length

$$c(n)\left(\frac{\sqrt{n}a}{r}\right)^{k+1}$$
.

There are at most  $r^n$  such cubes with points in  $D_k$ . The summed up volumes of the images of these cubes in  $\mathbb{R}^p$  are at most

$$c(n)^p \left(\frac{\sqrt{n}a}{r}\right)^{p(k+1)} r^n = cr^{n-p(k+1)}.$$

Since n - p(k+1) < 0, this will get arbitrary small for  $r \to \infty$ .

**Corollary B.13** (Brown). Let M and N be (finite dimensional) manifolds. Let  $f: M \to N$  be a differentiable  $(C^{\infty}$ -)maps. Then all the regular values of f lay dense in N.

We want to derive Brouwer's fixed point theorem from Sard's theorem.

**Definition B.14.** Let  $A \subset B$ . A retraction is a continuous map  $f : B \to A$ , so that  $f|_A = id$ , that is, f(x) = x for all  $x \in A$ .

**Theorem B.15.** There exists no retraction of  $\overline{B_1(0)} \subset \mathbb{R}^n$  on  $\mathbb{S}^{n-1}$ .

*Proof.* We prove the claim by contradiction. Let  $f: \overline{B_1(0)} \to \mathbb{S}^{n-1}$  be a retraction. Show at first, that then there also exists a  $C^{\infty}$ -retraction of  $\overline{B_1(0)}$  on  $\mathbb{S}^{n-1}$ : We find a retraction q, that is close to  $\partial B_1(0)$  of the class  $C^{\infty}$ , e.g.,

$$g(x) = \begin{cases} f\left(\frac{x}{|x|}\right) & \text{for } \frac{1}{2} \le |x| \le 1\\ f(2x) & \text{for } 0 \le |x| \le \frac{1}{2} \end{cases}.$$

Mollification in the interior gives a  $C^{\infty}$ -retraction. Hence we may assume that  $f \in C^{\infty}(\overline{B_1(0)}, \mathbb{S}^{n-1})$ . By Corollary B.13 there exists a regular value  $y \in \mathbb{S}^{n-1}$  of f. Hence the compact set  $f^{-1}(y)$  is a one-dimensional submanifold (first in  $B_1(0)$ , but since we can mollify f, also up to the boundary, since f is after construction constant on radial line segments close to  $\mathbb{S}^{n-1}$ ). Hence  $f^{-1}(y)$  is a one-dimensional manifold with boundary in  $\overline{B_1(0)}$ , whose boundary is a subset of  $\mathbb{S}^{n-1} = \partial B_1$ . It holds that  $y \in f^{-1}(y)$ , since f is a retraction. Let V be the component of  $f^{-1}(y)$  that contains y. Then V is a one-dimensional compact connected manifold and thus diffeomorph to a closed interval. Then y is the one boundary point of V. Let z be the other, which as well lays on  $\partial B_1(0)$ . It follows that z = f(z) in contradiction to  $y, z \in f^{-1}(y)$ .

**Theorem B.16** (Brouwer's fixed point theorem). Let  $f: \overline{B_1(0)} \to \overline{B_1(0)}$  be continuous. Then f has one fixed point, that is, there exists  $x \in \overline{B_1(0)}$  with f(x) = x.

*Proof.* If  $f(x) \neq x$  for all  $x \in \overline{B_1(0)}$ , we define g(x) to be the intersection of a line with  $\mathbb{S}^{n-1}$  beginning in f(x) through x. As constructed g is a retraction of  $\overline{B_1(0)}$  on  $\mathbb{S}^{n-1}$ .

# APPENDIX C. MAXIMUM PRINCIPLES

**Theorem C.1** (Strong elliptic maximum principle). Let M be closed and  $f: M \to \mathbb{R}$  satisfy

$$-\Delta_M f + b^i \nabla_i^M f + cf \le 0$$

for some smooth funtions  $b^i$  and  $c \le 0$ . If  $f \le 0$ , but  $f \not\equiv 0$ , then f < 0.

*Proof.* For a proof see [Eva02, §6.4, Theorem 4] or [Sch17b, Theorem 5.5] for  $M^n = \mathbb{R}^n$ .

Let  $M^n$  be a smooth n-dimensional manifold with boundary whose closure is compact. Let  $X: \overline{M}^n \times [0,T) \to \mathbb{R}^{n+m}$  be a family of smooth embeddings and set  $M_t := X(M^n,t)$ . For  $f \in C^{2;1}(M^n \times [0,T))$ , we define the parabolic operator

$$L(f) := \partial_t f - a^{ij} \nabla_i \nabla_i f - b^i \nabla_i f - cf, \qquad (35)$$

where  $a_{ij}, b_i, c \in L^{\infty}$  may depend on  $p, t, (g_{kl})_{kl}, f, \nabla f$ , and  $\nabla^2 f$ , and where  $(a^{ij})_{ij}$  is positive semi-definite, that is,

$$\lambda |\xi|^2 \le a_{ij} \xi_i \xi_j \le \Lambda |\xi|^2 \tag{36}$$

for all  $\xi \in \mathbb{R}^n$ . For R > 0,  $p_0 \in M^n$  and  $t_0 \in [0, T)$ , define the spatial neighbourhood

$$U_R(p_0, t_0) := X^{-1} (B_R(X(p_0, t_0)) \cap M_{t_0})$$
  
=  $\{ p \in M^n \mid |X(p, t_0) - X(p_0, t_0)| < R \},$ 

the parabolic neighbourhood

$$Q_R(p_0, t_0) := \left\{ (p, t) \in M^n \times \left( t_0 - R^2, t_0 \right] \mid |X(p, t) - X(p_0, t)| < R \right\}$$
$$= \bigcup_{t \in (t_0 - R^2, t_0]} (U_R(p_0, t) \times \{t\})$$

and, for an open set  $U \subset M^n$  and  $[t_1, t_0] \subset [0, T)$ , the parabolic boundary

$$\mathcal{P}(U \times [t_1, t_0]) := (U \times \{t_1\}) \cup (\partial U \times (t_1, t_0]).$$

**Theorem C.2** (Weak parabolic maximum principle). Let  $U \subset M^n$  be open and let  $f \in C^{2;1}(Q) \cap C^0(\mathcal{P}Q)$  for  $Q := U \times [t_1, t_0]$ .

- (i) If L(f) > 0 on Q and f > 0 on  $\mathcal{P}Q$ . Then f > 0 in Q.
- (ii) If  $L(f) \leq 0$  on Q and  $f \leq 0$  on  $\mathcal{P}Q$ . Then  $f \leq 0$  in Q.

**Theorem C.3** (Strong parabolic maximum principle). Let  $U \subset M^n$  be open,  $Q := U \times [0,T)$ , and  $f \in C^{2;1}(Q) \cap C^0(\bar{Q})$ .

- (i) Let  $L(f) \ge 0$  in Q. If there exists  $(p_0, t_0) \in Q \setminus \mathcal{P}Q$  with  $f(p_0, t_0) = \min_{\bar{Q}} f$ , then f is constant in  $\bar{Q}$ .
- (ii) Let  $L(f) \leq 0$  in Q. If there exists  $(p_0, t_0) \in Q \setminus PQ$  with  $f(p_0, t_0) = \max_{\bar{Q}} f$ , then f is constant in  $\bar{Q}$ .

C.1. **2-tensors.** We follow the lines of [CCG<sup>+</sup>08, Chapter 12]. Let T > 0 and  $(M^n, g(t))_{t \in [0,T)}$  a closed manifold with a family of metrics, that depend smoothly on time. Let  $m = (m_{ij})_{1 \le i,j \le n}$  be symmetric with  $m_{ij} \in C^{\infty}(M^n \times [0,T))$ . Let  $b = (b_{ij}(m,p,t))_{1 \le i,j \le n}$  be symmetric with  $b_{ij} \in C^1(M^n \times [0,T))$  and satisfy the null eigenvector condition, that is, if  $m_{ij}\xi^j = 0$  for  $1 \le i \le n$  then also  $b_{ij}\xi^i\xi^j \ge 0$ . Let  $u^k \in L^{\infty}(M^n \times [0,T))$ ,  $1 \le k \le n$ .

Theorem C.4 (Weak parabolic maximum principle for 2-tensors). Let

$$\partial_t m_{ij} \succeq \Delta_{g(t)} m_{ij} + u^k \nabla_k^{g(t)} m_{ij} + b_{ij} (m_{kl}, \cdot)$$

in  $M^n \times (0,T)$  and  $m_{ij}(\cdot,0) \succeq 0$ . Then  $m_{ij}(\cdot,t) \succeq 0$  for  $0 \leq t < T$ .

**Theorem C.5** (Strong parabolic maximum principle for 2-tensors I). Let b be locally Lipschitz in m. Let

$$\partial_t m_{ij} = \Delta_{g(t)} m_{ij} + u^k \nabla_k^{g(t)} m_{ij} + b_{ij} (m_{kl}, \cdot)$$

in  $M^n \times (0,T)$ ,  $m_{ij}(\cdot,0) \succeq 0$  for all  $t \in [0,T)$  and  $m_{ij}(p_0,0) \succ 0$  for  $p_0 \in M^n$ . Then  $m_{ij}(\cdot,t) \succ 0$  for 0 < t < T.

*Proof.* Let  $p \in M^n$  and  $U \subset M^n$  so that  $p, p_0 \in U$  and so that  $\overline{U}$  is a compact manifold with smooth boundary. Define  $\varphi_1 : \overline{U} \times [0,T) \to \mathbb{R}$  by

$$\varphi_1 \le \lambda_1(\cdot, 0) \quad \text{in } \overline{U}$$

$$\varphi_1 \equiv 0 \quad \text{on } \partial U$$

$$2\varphi_1(p_0) \ge \lambda_1(p_0, 0).$$

Let C>0 to be chosen later and let  $f:\overline{U}\times[0,T)\to\mathbb{R}$  a solution of

$$\begin{split} \partial_t f &= \Delta_{g(t)} f + u^k \nabla_k^{g(t)} f - C f & \text{in } U \times (0, T) \\ f &\equiv 0 & \text{on } \partial U \times [0, T) \\ f(\cdot, 0) &= \varphi_1 & \text{in } U \,. \end{split}$$

Since  $m_{ij}(p_0,0) > 0$ , we also have  $\varphi_1(p_0) > 0$ . The strong maximum principle for functions, Theorem C.3, yields that f > 0 in  $U \times (0,T)$ . The weak maximum principle, Theorem C.2, yields

$$f(x,t) \le \max_{p \in \overline{U}} \varphi_1(x) \le \max_{p \in \overline{U}} \lambda_1(x,0)$$

in  $U \times (0,T)$ . Define the tensor

$$\tilde{m}_{ij} = m_{ij} + (\varepsilon e^{Ct} - f)\delta_{ij} \,,$$

where  $\varepsilon > 0$ . Then

$$\tilde{m}_{ij} \succeq \lambda_1 \delta_{ij} + (\varepsilon e^{Ct} - \lambda_1) \delta_{ij} \succ 0$$

and

$$\partial_{t}\tilde{m}_{ij} = \partial_{t}m_{ij} + \left(\varepsilon C e^{Ct} - \partial_{t}f\right)\delta_{ij}$$

$$= \Delta_{g(t)}(m_{ij} - f\delta_{ij}) + u^{k}\nabla_{k}^{g(t)}(m_{ij} - f\delta_{ij})$$

$$+ b_{ij}(m_{kl}) + C\left(\varepsilon e^{Ct} + f\right)\delta_{ij}$$

$$= \Delta_{g(t)}\tilde{m}_{ij} + u^{k}\nabla_{k}^{g(t)}\tilde{m}_{ij} + b_{ij}(\tilde{m}_{kl})$$

$$- \left(b_{ij}(\tilde{m}_{kl}) - b_{ij}(m_{kl})\right) + C\left(\varepsilon e^{Ct} + f\right)\delta_{ij}.$$

Since  $b_{ij}$  is Lipschitz in  $m_{ij}$ ,

$$b_{ij}(\tilde{m}_{kl}) - b_{ij}(m_{kl}) \leq \operatorname{Lip}(b_{kl})(\tilde{m}_{ij} - m_{ij}) = \operatorname{Lip}(b_{kl})(\varepsilon e^{Ct} + f)\delta_{ij}$$
.

By choosing  $C \geq \text{Lip}(b_{ij})$  and  $\varepsilon$  such that  $\varepsilon \leq e^{-Ct}$ , we obtain

$$\partial_t \tilde{m}_{ij} \succeq \Delta_{g(t)} \tilde{m}_{ij} + u^k \nabla_k^{g(t)} \tilde{m}_{ij} + b_{ij} (\tilde{m}_{kl})$$

$$+ (C - \operatorname{Lip}(b_{ij})) \left( \varepsilon e^{Ct} + f \right) \delta_{ij}$$

$$\succeq \Delta_{g(t)} \tilde{m}_{ij} + u^k \nabla_k^{g(t)} \tilde{m}_{ij} + b_{ij} (\tilde{m}_{kl}).$$

The weak maximum principle, Theorem C.4, implies  $\tilde{m}_{ij} \succeq 0$  on  $\overline{U} \times [0,T)$  for  $\varepsilon \in (0, e^{-Ctt}]$ . Thus  $m_{ij} \succeq (-\varepsilon e^{Ct} + f)\delta_{ij}$  on  $\overline{U} \times [0,T)$  for  $\varepsilon \in (0, e^{-Ctt}]$ . Letting  $\varepsilon \to 0$  yields  $m_{ij} \succeq f\delta_{ij} \succ 0$  on  $\overline{U} \times [0,T)$ .

Theorem C.6 (Strong parabolic maximum principle for 2-tensors II). Let

$$\phi_k(p,t) := \inf_{\{\tau_1,\dots,\tau_k\} \text{ orthonormal}} (m(\tau_1,\tau_1) + \dots + m(\tau_k,\tau_k))$$
$$= \lambda_1(p,t) + \dots + \lambda_k(p,t)$$

where  $k \in \{1, ..., n\}$ . Let b be locally Lipschitz in m. Let

$$\partial_t m_{ij} = \Delta_{g(t)} m_{ij} + u^k \nabla_k^{g(t)} m_{ij} + b_{ij} (m_{kl}, \cdot)$$

in  $M^n \times (0,T)$ ,  $\phi_k(\cdot,0) \ge 0$  in  $M^n$  and  $\phi_k(p_0,0) > 0$  for  $k \in \{1,\ldots,n\}$  and  $p_0 \in M^n$ . Then  $\phi_k(\cdot,t) > 0$  for 0 < t < T.

*Proof.* Let  $p \in M^n$  and  $U \subset M^n$  so that  $p, p_0 \in U$  and so that  $\overline{U}$  is a compact manifold with smooth boundary. Define  $\varphi_k : \overline{U} \times [0,T) \to \mathbb{R}$  by

$$k\varphi_k \le \phi_k(\cdot, 0)$$
 in  $\overline{U}$   
 $\varphi_k \equiv 0$  on  $\partial U$   
 $k\varphi_k(p_0) \ge \lambda_1(p_0, 0)$ .

Let C>0 to be chosen later and let  $f:\overline{U}\times[0,T)\to\mathbb{R}$  a solution of

$$\begin{split} \partial_t f &= \Delta_{g(t)} f + u^k \nabla_k^{g(t)} f - C f &\quad \text{in } U \times (0, T) \\ f &\equiv 0 &\quad \text{on } \partial U \times [0, T) \\ f(\cdot, 0) &= \varphi_k &\quad \text{in } U \,. \end{split}$$

Since  $\phi_k(p_0,0) > 0$ , we also have  $\varphi_k(p_0) > 0$ . The strong maximum principle for functions, Theorem C.3, yields that f > 0 in  $U \times (0,T)$ . The weak maximum principle, Theorem C.2, yields

$$f(x,t) \le \max_{p \in \overline{U}} \varphi_k(x) \le \max_{p \in \overline{U}} \phi_k(x,0)$$

in  $U \times (0,T)$ . Define the tensor

$$\tilde{m}_{ij} = m_{ij} + (\varepsilon e^{Ct} - f)\delta_{ij} \,,$$

for  $\varepsilon > 0$  and

$$\tilde{\phi}_k(p,t) := \inf_{\{\tau_1,\dots,\tau_k\} \text{ orthonormal}} \left( \tilde{m}(\tau_1,\tau_1) + \dots + \tilde{m}(\tau_k,\tau_k) \right)$$
$$= \phi_k(x,t) + k(\varepsilon e^{Ct} - f(x,t)).$$

We want to show that  $\tilde{\phi}_k > 0$  on  $\overline{U} \times [0,T)$  for  $\varepsilon > 0$  small enough. Assume the opposite. Since  $\tilde{\phi}_k > 0$  in  $U \times \{0\}$  and  $\partial U \times [0,T)$ , there exists a point  $(p_1,t_1) \in U \times [0,T)$  with

$$\tilde{\phi}_k(p_1.t_1) = 0$$
 and  $\tilde{\phi}_k(p.t) > 0$  for all  $(p,t) \in U \times [0,t_1)$ .

Let  $\boldsymbol{\tau}_1^0, \dots \boldsymbol{\tau}_k^0 \in T_{p_1} M^n$  be orthonormal with

$$\tilde{m}(\boldsymbol{\tau}_1^0, \boldsymbol{\tau}_1^0) + \dots + \tilde{m}(\boldsymbol{\tau}_k^0, \boldsymbol{\tau}_k^0) = 0$$

in  $(p_1, t_1)$ . Extend each  $\tau_i^0$  in space and time to a lokal vector field  $\tau_i$  by parallel translation of  $\tau_i^0$  along geodesics starting from  $p_1$  with respect to  $\nabla^{g(t_1)}$  and constant in time. Then

$$\nabla \boldsymbol{\tau}_i(p_1, t_1) = 0$$
,  $\Delta \boldsymbol{\tau}_i(p_1, t_1) = 0$ ,  $\partial_t \boldsymbol{\tau}_i(p_1, t_1) = 0$ .

Define in a neighbourhood of  $(p_1, t_1)$ 

$$\psi_k(p,t) := \tilde{m}(p,t)(\tau_1,\tau_1) + \dots + \tilde{m}(p,t)(\tau_k,\tau_k)$$

where  $\psi_k(p_1, t_1) = 0$  and

$$\psi_k(p,t) \ge \tilde{\phi}_k(p,t) \ge 0$$

for all  $p \in U$  and  $t \in [0, t_1]$ . At  $(p_1, t_1)$ , we have

$$0 \geq (\partial_t - \Delta - u^l \nabla_l) \psi_k$$

$$= \sum_{i=1}^k (\partial_t - \Delta - u^l \nabla_l) \tilde{m}(\boldsymbol{\tau}_i^0, \boldsymbol{\tau}_i^0)$$

$$= \sum_{i=1}^k b(\tilde{m})(\boldsymbol{\tau}_i^0, \boldsymbol{\tau}_i^0) - \sum_{i=1}^k (b(\tilde{m}) - b(m)) (\boldsymbol{\tau}_i^0, \boldsymbol{\tau}_i^0) + C \left(\varepsilon e^{Ct} + f\right)$$

$$\geq \left(kC - \sum_{i=0}^k \operatorname{Lip}(b)(\boldsymbol{\tau}_i^0, \boldsymbol{\tau}_i^0)\right) \left(\varepsilon e^{Ct} + f\right) > 0$$

if we choose  $C \geq \text{Lip}(b_{ij})$  and  $\varepsilon$  such that  $\varepsilon \leq e^{-Ct}$ . This is a contradiction. Hence,  $\tilde{\phi}_k > 0$  on  $\overline{U} \times [0,T)$  for  $\varepsilon \leq e^{-Ct}$ . Thus  $\phi_k \geq -k(\varepsilon e^{Ct} - f)$  on  $\overline{U} \times [0,T)$  for  $\varepsilon \in (0, e^{-Ctt}]$ . Letting  $\varepsilon \to 0$  yields  $\psi_k \geq f > 0$  on  $\overline{U} \times [0,T)$ .

**Theorem C.7** (Strong parabolic maximum principle for 2-tensors III). Let b be locally Lipschitz in m. Let

$$\partial_t m_{ij} = \Delta_{g(t)} m_{ij} + u^k \nabla_k^{g(t)} m_{ij} + b_{ij} (m_{kl}, \cdot)$$

in  $M^n \times (0,T)$  and  $m_{ij}(\cdot,0) \succeq 0$  for all  $t \in [0,T)$ . Then

(i) If  $t_2 > t_1$  in [0, T), then

$$\inf_{p \in M^n} \operatorname{rank} m(p, t_2) \ge \sup_{p \in M^n} \operatorname{rank} m(p, t_1)$$

and there exists  $\delta > 0$  so that rank m(p,t) is constant for all  $p \in M^n$  and  $t \in (0,\delta)$ .

- (ii) ( ker m is smooth in space and time). Let  $(0, \delta)$  be the time interval from (i). Then, ker  $m(t) \subset TM^n$  is a smooth subspace which depends smoothly on time for  $t \in (0, \delta)$ .
- (iii) (ker m is parallel in space and time). Let  $(0, \delta)$  be the time interval from (i). Then, ker m(t) is invariant under parallel transport in space and constant in time for  $t \in (0, \delta)$ .

Proof.  $\Box$ 

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