

MEAN CURVATURE FLOW

FRIEDERIKE DITTBERNER

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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 $X_0 : M^n \rightarrow \mathbb{R}^{n+1}$, where M^n is an abstract smooth manifold. We consider the family of embeddings $X : M^n \times [0, T) \rightarrow \mathbb{R}^{n+1}$ with

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

for all $p \in M^n$ and

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

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 $M_t = \mathbb{S}_{r(t)}^{n-m} \times \mathbb{R}^m$ with $r(t) = \sqrt{R^2 - 2(n-m)t}$ exist for $t \in (-\infty, R^2/2(n-m))$.

(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

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

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 $\tilde{X}(\cdot, t) : M^n \rightarrow \mathbb{R}^{n+1}$ with $M_t = \tilde{X}(M^n, t)$ be a family of embeddings satisfying the equation

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

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

$$D^{\mathbb{R}^{n+1}} \tilde{X}(\phi(p, t), t) \partial_t \phi(p, t) = - \left(\partial_t \tilde{X}(\phi(p, t), t) \right)^\top,$$

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

$$\partial_t X(p, t) = \partial_t \tilde{X}(p, t) + D^{\mathbb{R}^{n+1}} \tilde{X}(\phi(p, t), t) \partial_t \phi(p, t) = \left(\partial_t \tilde{X}(q, t) \right)^\perp = \mathbf{H}(X(p, t)).$$

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

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

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

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,*

$$\begin{aligned} \partial_t g_{ij} &= -H h_{ij}, \\ \partial_t d\mu_t^n &= -H^2 d\mu_t^n, \\ \partial_t \nu &= \nabla H, \\ \partial_t h_{ij} &= \nabla_i \nabla_j H - H g^{km} h_{ik} h_{jm}, \\ \partial_t H &= \Delta H + H |A|^2, \\ \partial_t |A|^2 &= \Delta |A|^2 - |\nabla A|^2 + 2|A|^4, \\ \partial_t |\nabla^m A|^2 &\leq \Delta |\nabla^m A|^2 - 2|\nabla^{m+1} A|^2 \\ &\quad + C(m, n) \sum_{i+j+k=m} |\nabla^m A| \cdot |\nabla^i A| \cdot |\nabla^j A| \cdot |\nabla^k A| \end{aligned}$$

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

Proof. See e.g. [Sch18, Section 3]. □

Corollary 1.5. *The mean curvature flow 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 $\partial_t g_{ij} = -2Fh_{ij}$ and

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

with equality if and only if $F = H$. \square

Theorem 1.6 (Comparison principle). *Let $X : M^n \times [0, T) \rightarrow \mathbb{R}^{n+1}$ and $Y : N^n \times [0, T) \rightarrow \mathbb{R}^{n+1}$ 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 $v = v_1 \oplus v_2 \in T_{X(p_t, t)} M_t \oplus T_{Y(q_t, t)} N_t$ 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 $f, h : U \times (t - \varepsilon, t + \varepsilon) \rightarrow \mathbb{R}$, where $U \subset \mathbb{R}^n$. After rotation, we can assume that $\text{span}(e_1, \dots, e_n) \subset \mathbb{R}^{n+1}$ 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

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

and

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

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

$$\Delta f(0, t) - \Delta h(0, t) \geq 0 \quad \text{and} \quad Df(0, t) = Dh(0, t) = 0$$

and so

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

Since

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

we obtain at (p_t, q_t) that

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

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

Proposition 1.7 (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

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

and

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

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

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

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 $t_i \rightarrow 0$. We have for a subsequence, that $p_i \rightarrow p$ and $q_i \rightarrow q$. 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 $m^* : [0, T] \rightarrow \mathbb{R}_+$ such that

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

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

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

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

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

and the geodesic neighbourhood

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

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

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

Consider the open set

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

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 1.6, shows that $\partial_t d(t) \geq 0$. Hence

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

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

Theorem 1.8 (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. $H \geq 0$. Then $H > 0$ for all $t \in (0, T)$.*

Proof. See [Sch17d, 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.5. 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 $\lambda_1 > 0$. \square

2. HOMOTHETICALLY SHRINKING SOLUTIONS

Definition 2.1 (Homothetically shrinking solutions, Brakke [Bra78, Appendix C]). Let $\lambda : [t_0, T] \rightarrow \mathbb{R}_+$ be smooth and decreasing, $\lambda(t_0) = 1$ and $\lambda(T) = 0$. Let $x_0 \in \mathbb{R}^{n+1}$. A *homothetically shrinking* solution $X : M^n \times [t_0, T] \rightarrow \mathbb{R}^{n+1}$ 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 $\tilde{X} : M^n \times [t_0, T] \rightarrow \mathbb{R}^{n+1}$ with $M_t = \tilde{X}(M^n, t)$ satisfying the evolution equation

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

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

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

for $q \in M^n$, where the embeddings $X(\cdot, t) : M^n \rightarrow \mathbb{R}^{n+1}$ 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 $X(p, t) = \sqrt{-2t}X(p, t_0)$. \square

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 $H(x) = \langle x, \boldsymbol{\nu} \rangle / 2$ 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 $\mathbb{S}_{\sqrt{2m}}^m \times \mathbb{R}^{n-m}$ for $m = 0, 1, \dots, n$.*

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 $\nabla_k h_{ij} = \nabla_j h_{ik}$ at one fixed point where the Christoffel symbols vanish, that

$$\begin{aligned} \nabla_i \nabla_j H &= g_{ik} h_j^k + \langle x, \nabla_i \nabla_k x \rangle h_j^k + \langle x, \partial_k K \rangle \nabla_i h_j^k \\ &= h_{ij} + \langle x, \nu \rangle h_{ik} h_j^k + \langle x, \partial_k x \rangle g^{kl} \nabla_i h_{jl} \\ &= h_{ij} - H h_{ik} h_j^k + \langle x, \partial_k x \rangle g^{kl} \nabla_l h_{ij} \\ &= h_{ij} - H h_{ik} h_j^k + \langle x, \nabla h_{ij} \rangle. \end{aligned} \quad (1)$$

Contracting with g^{ij} we have

$$\Delta H = H (1 - |A|^2) + \langle x, \nabla H \rangle. \quad (2)$$

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

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

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

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

Assume that $H = 0$. As M is complete and x is a tangent vector field on M by the equation $\langle x, \nu \rangle = 0$, 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

$$\eta_R = \nu_{\partial(M \cap B_R(0))}$$

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

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

and similarly

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

Hence,

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

As we have

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

the quantity σ_R is nonnegative and nondecreasing in R . If now we show that

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

we can conclude that, at every point of M ,

$$4|A|^2|\nabla A|^2 = |\nabla |A|^2|^2. \quad (3)$$

Getting back to the definitions of ε_R and δ_R , we have

$$\begin{aligned}
|\sigma_R| &= \left| -2 \int_{\partial(M \cap B_R(0))} \frac{H}{|A|} \langle \nabla |A|^2, \eta \rangle \exp\left(-\frac{R^2}{2}\right) d\mu^{n-1} \right. \\
&\quad \left. + 4 \int_{\partial(M \cap B_R(0))} |A| \langle \nabla H, \eta \rangle \exp\left(-\frac{R^2}{2}\right) d\mu^{n-1} \right| \\
&\leq 4 \exp\left(-\frac{R^2}{2}\right) \int_{\partial(M \cap B_R(0))} \left(\frac{H}{|A|} |\nabla |A|^2| + |A| |\nabla H| \right) d\mu^{n-1} \\
&\leq 8 \exp\left(-\frac{R^2}{2}\right) \int_{\partial(M \cap B_R(0))} (H |\nabla A| + |A| |\nabla H|) d\mu^{n-1} \\
&\leq C \exp\left(-\frac{R^2}{2}\right) \mu^{n-1}(\partial(M \cap B_R(0))),
\end{aligned}$$

by the estimates on A and ∇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

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

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 $R_0 < r < R$,

$$\begin{aligned}
\mu^n(M \cap B_R(0)) - \mu^n(M \cap B_r(0)) &= \int_r^R \frac{d}{d\xi} \mu^n(M \cap B_\xi(0)) d\xi \\
&= \int_r^R \int_{M \cap B_\xi(0)} \operatorname{div}_{M \cap B_\xi(0)} \eta_\xi d\mu^{n-1} d\xi \\
&= - \int_r^R \int_{M \cap B_\xi(0)} \langle \eta_\xi, \mathbf{H}_{M \cap B_\xi(0)} \rangle d\mu^{n-1} d\xi \\
&\quad + \int_r^R \int_{\partial(M \cap B_\xi(0))} \langle \eta_\xi, \eta_\xi \rangle d\mu^{n-1} d\xi \\
&= \int_r^R \int_{\partial(M \cap B_\xi(0))} d\mu^{n-1} d\xi \\
&\geq \delta \int_r^R \xi \exp\left(\frac{\xi^2}{4}\right) d\xi = 2\delta \left(\exp\left(\frac{R^2}{4}\right) - \exp\left(\frac{r^2}{4}\right) \right).
\end{aligned}$$

Then

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

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

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

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

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

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 ∇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

$$\nabla_k H = c_k H \quad \text{and} \quad \nabla_k |A|^2 = 2c_k |A|^2,$$

hence

$$\nabla_k \log H = c_k \quad \text{and} \quad \nabla_k \log |A|^2 = 2c_k.$$

This implies

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

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

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

at one point, where the Christoffel symbols vanish. Computing then in normal coordinates with an orthonormal basis $\{\tau_1, \dots, \tau_n\}$ such that $\tau_1 = \nabla H / |\nabla H|$, we have with $g^{ij} = \delta^{ij}$,

$$\begin{aligned} 0 &= |h_{ij} \nabla_k H - h_{ik} \nabla_j H|^2 \\ &= (h_{ij} \nabla_k H - h_{ik} \nabla_j H) g^{il} g^{jm} g^{kn} (h_{lm} \nabla_n H - h_{ln} \nabla_m H) \\ &= 2|\nabla H|^2 |A|^2 - 2g^{il} g^{jm} g^{kn} h_{ij} h_{ln} \nabla_k H \nabla_m H \\ &= 2|\nabla H|^2 |A|^2 - 2g^{il} h_i^m h_l^k \nabla_k H \nabla_m H \\ &= 2|\nabla H|^2 |A|^2 - 2g^{il} h_i^1 h_l^1 \nabla_1 H \nabla_1 H \\ &= 2|\nabla H|^2 \left(|A|^2 - \sum_{i=1}^n (h_i^1)^2 \right). \end{aligned}$$

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 λ_m with eigenvector ξ_m ,

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

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

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

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

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

by means of the structural equation $H = \langle x, \nu \rangle$. 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 $M = \mathbb{S}_{\sqrt{n}}^n$.

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 $k = 1, \dots, n - m$. 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 Gronwall's 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 γ 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 κ be the curvature of γ in \mathbb{R}^{n+1} , we have

$$\kappa = -\langle \nu_M, \partial_s^2 \gamma \rangle = h_{ij} \partial_s \gamma^i \partial_s \gamma^j = 0,$$

thus γ 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 $v_k(x)$, $k = 1, \dots, n - m$, along σ , 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 Gronwall's lemma it follows that $h_{ij}(\gamma(s))v_k^j(\gamma(s)) = 0$ for every $s \in \mathbb{R}$ and $k = 1, \dots, n - m$, 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 σ , we have

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

so the extended vectors v_k are constant in \mathbb{R}^{n+1} , which means that the parallel extension is independent of the geodesic σ , 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 $f : M \rightarrow \ker A$ by

$$f(x) = v.$$

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

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

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 ν_M to M at a point of L stays in $(\ker A)^\perp$ so it must coincide with the normal ν_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 $L = \mathbb{S}_{\sqrt{m}}^m$ and $M = \mathbb{S}_{\sqrt{m}}^m \times \mathbb{R}^{n-m}$ which proves the claim. \square

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 $\mathbb{S}_{\sqrt{2m}}^m \times \mathbb{R}^{n-m}$ for $m = 0, 1, \dots, n$.*

2.2. Curves.

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

Proof. See [Man11, Proposition 3.4.1]. We scale the curve by the factor $1/2$ so that $\kappa = \langle x, \nu \rangle$ 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 $\tau = \partial_s X$ and a unit normal vector field $\nu = (\tau_2, -\tau_1)$, which is the clockwise rotation of $\pi/2$ in \mathbb{R}^2 of the vector τ . Then the curvature is given by

$$\kappa = -\langle \partial_s \tau, \nu \rangle = \langle \tau, \partial_s \nu \rangle$$

so that

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

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

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

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 κ is identically zero so that Σ is a line. Since $\langle x, \nu \rangle = 0$ for every $x \in \Sigma$, we conclude that $0 \in \Sigma$. So we suppose that κ 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 $|X|^2$,

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

we get

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

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

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

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

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

so that,

$$(\partial_\vartheta \kappa)^2 + \kappa^2 - \log \kappa^2 \equiv E \geq 1$$

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 κ 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, Σ must be closed, simple and strictly convex, as $\kappa > 0$ at every point.

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

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

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

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

and

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

Furthermore, integration by parts yields

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

and

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

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

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

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

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

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

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

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

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

Then $|I_i| \leq \pi/2$ for all $i \in \{1, \dots, i_0\}$. 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 κ would imply $\partial_{\vartheta} \kappa = 0$ everywhere, hence $\kappa = 1$ identically and we would be in the case of the unit circle. Suppose that Σ 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 κ 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

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

so that

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

We compute

$$\begin{aligned} 4 \int_0^{\vartheta_0} \sin(2\vartheta) \frac{\partial_{\vartheta} \kappa}{\kappa} d\vartheta &= \int_0^{\vartheta_0} \sin(2\vartheta) (\partial_{\vartheta}^3 \kappa^2 + 4\partial_{\vartheta} \kappa^2) d\vartheta \\ &= \sin(2\vartheta) \partial_{\vartheta}^2 \kappa^2 \Big|_0^{\vartheta_0} - 2 \int_0^{\vartheta_0} \cos(2\vartheta) \partial_{\vartheta}^2 \kappa^2 d\vartheta + 4 \int_0^{\vartheta_0} \sin(2\vartheta) \partial_{\vartheta} \kappa^2 d\vartheta \\ &= 2 \sin(2\vartheta_0) (\kappa(\vartheta_0) \partial_{\vartheta}^2 \kappa(\vartheta_0) + (\partial_{\vartheta} \kappa)^2(\vartheta_0)) - 2 \cos(2\vartheta) \partial_{\vartheta} \kappa^2 \Big|_0^{\vartheta_0} \\ &\quad - 4 \int_0^{\vartheta_0} \sin(2\vartheta) \partial_{\vartheta} \kappa^2 d\vartheta + 4 \int_0^{\vartheta_0} \sin(2\vartheta) \partial_{\vartheta} \kappa^2 d\vartheta \\ &= 2 \sin(2\vartheta_0) (\kappa(\vartheta_0) \partial_{\vartheta}^2 \kappa(\vartheta_0) + (\partial_{\vartheta} \kappa)^2(\vartheta_0)) \\ &\quad - 4 \cos(2\vartheta_0) \kappa(\vartheta_0) \partial_{\vartheta} \kappa(\vartheta_0) + 4\kappa(0) \partial_{\vartheta} \kappa(0). \end{aligned}$$

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 $0 < 2\vartheta_0 \leq \pi$. Looking at the left-hand integral we see instead that the factor $\sin(2\vartheta)$ is always nonnegative, since $2\vartheta_0 \leq \pi$ 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 κ at 0 to a local minimum of κ at ϑ_0 without crossing any other critical point of κ . This gives a contradiction so Σ must be the unit circle. \square

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 $X, Y : M^n \rightarrow \mathbb{R}^{n+1}$ are *conformal*, if there exists $\lambda : M^n \rightarrow \mathbb{R}$ with

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

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

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

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 open set $U \subset \mathbb{R}^n$ of a strictly convex function

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

so that $y \rightarrow \infty$ as $x = (x_1, \dots, x_n) \rightarrow \partial U$. By translating upwards if necessary, since y is bounded below, we can assume $y \geq e$ everywhere, so that $\log \log y \geq 0$. 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 $\gamma : I \rightarrow M$ going to infinity, we have $\gamma^n \rightarrow \infty$. Therefore its length satisfies,

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

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

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

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

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

and let dV_y denote the volume element of the part of M between height y and height $y + dy$. We can divide M into pieces M_1, \dots, M_n , where $\frac{\partial y}{\partial x^i}$ is largest on M_i , and estimate dV_y^i from above on each piece. On M_i , 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 $\left| \frac{\partial y}{\partial x^i} \right| \geq \delta > 0$,

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

and thus

$$dV_y^i \leq C dx^1 \dots dx^{i-1} dy dx^{i+1} \dots dx^n$$

on M_i . By the gradient estimate shows that

$$|x| \leq Cy$$

for a suitable large constant. Since the region of integration over $x^1, \dots, x^{i-1}, x^{i+1}, \dots, x^n$ to get dV_y from dV has its volume bounded by Cy^{n-1} ,

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

Hence,

$$d\tilde{V}_y^i \leq \frac{C dy}{y \log^n y}$$

and

$$\tilde{V} \leq C \int_e^\infty \frac{dy}{y \log^n y} = \frac{-C}{(n-1) \log^{n-1} y} \Big|_e^\infty = \frac{C}{n-1} < \infty. \quad \square$$

Remark 3.4. 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 $\vartheta \in [-\pi/2 + \varepsilon, \pi/2 - \varepsilon]$ and $z(\vartheta) \in [-R, R]$, where we choose $\varepsilon, R(\varepsilon) > 0$. 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

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

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, simply connected 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 quasis-conformal 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 . We can assume p_N to be the north pole and p_S the south pole. (If no two such points can be found, the closure of U has to be the whole \mathbb{S}^n .) We can then find an $\varepsilon > 0$ so that the ε -ball around p_N lies in U ,

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

while the ε -ball around p_S avoids U ,

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

We can then find a conformal map φ of the sphere \mathbb{S}^n minus these two balls to the cylinder $\mathbb{S}^{n-1} \times [0, L]$,

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

taking the boundary of the ε -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 ε -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 ε -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

$$g^U \quad \text{on} \quad U$$

with finite volume conformal to the round metric

$$g^{\mathbb{S}^n} \quad \text{on} \quad \mathbb{S}^n$$

would give a complete metric with finite volume on

$$g^W \quad \text{on} \quad W$$

conformal to the product metric

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

We show that such cannot exist. We introduce coordinates

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

and

$$t \quad \text{on} \quad [0, L]$$

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 (ϑ, 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 conformally equivalent metric on W is given by

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

for some function λ 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

$$\iint_W \lambda d\mu dt \leq \left(\iint_W \lambda^n d\mu dt \right)^{1/n} \left(\iint_W d\mu dt \right)^{(n-1)/n}$$

and surely

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

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 ϑ is constant we have $\tilde{g} = \lambda$. Thus there is some ϑ 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. \square

Theorem 3.6 (Hamilton, [Ham94, Theorem 1.1]). *Let M be a smooth strictly convex, simply connected hypersurface bounding a region in \mathbb{R}^{n+1} , $n \geq 2$. Suppose that its second fundamental form is ε -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 $\nu : M \rightarrow \mathbb{S}^n$ gives a diffeomorphism of the convex hypersurface M^n 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 on U conformal to the round metric which is complete with finite volume. However, the pinching condition implies

$$\tilde{h}_i^k \geq \varepsilon H \delta_i^k$$

so that

$$D\nu(\tilde{\partial}_i) = \tilde{h}_i^k \tilde{\partial}_k \geq \varepsilon H \delta_i^k \tilde{\partial}_k = \varepsilon H \tilde{\partial}_i$$

and

$$\langle D\nu(\tilde{\partial}_i), D\nu(\tilde{\partial}_j) \rangle \geq (\varepsilon H)^2 \tilde{g}_{ij}$$

so that the Weingarten map $D\nu : TM \rightarrow T\mathbb{S}^n$, induces metric which is quasi-conformal with respect to \tilde{g} . This is a contradiction. \square

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 $X(p_k, t_k) \rightarrow x_0$ for $k \rightarrow \infty$.

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 $X(p_k, t_k) \rightarrow x$ for $k \rightarrow \infty$. 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 $|A|^2 \leq C_0$ on $M^n \times [0, T)$, then $|\nabla^m A|^2 \leq C_m$ on $M^n \times [0, T)$, where $C_m = C_m(n, M_0, C_0)$.*

Proof. See [Sch17d, 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

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

and

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

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

$$\begin{aligned} \partial_t f &\leq \Delta f - B_m (1 + |\nabla^m A|^2) - 2 B_m |\nabla^m A|^2 + B_{m-1} B_m \\ &\leq \Delta f - B_m f + B_m^2 |\nabla^{m-1} A|^2 + B_{m-1} B_m \\ &\leq \Delta f - B_m f + B. \end{aligned}$$

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

$$\begin{aligned} \partial_t \tilde{f} &\leq \exp(B_m t) (B_m f + \partial_t f) - \exp(B_m T) B \\ &\leq \exp(B_m t) (\Delta f + B) - \exp(B_m T) B \leq \Delta \tilde{f} \end{aligned}$$

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

$$f(\cdot, t) \leq \exp(-B_m t) \left(\max_M \tilde{f}(\cdot, 0) + \exp(B_m T) B t \right) \leq C. \quad \square$$

Theorem 4.3. *Let $T < \infty$ and $(M_t)_{t \in [0, T]}$ be a family of smooth, immersed hyper-surfaces 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

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

Then M_T is smooth.

Proof. By Proposition 4.2,

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

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

$$\partial_t \nu = \nabla H$$

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,

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

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 $f : U_{r,t_0}(p) \times [t_0, T) \rightarrow \mathbb{R}$ 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 Arzelà–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 $i \in \{1, \dots, N_0\}$. 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, $X_\infty^l = X^l(\cdot, T) : \bigcup_{i=1}^{N_l} U_{r,t_0}(p_i) \rightarrow \mathbb{R}^{n+1}$ 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 $X_\infty^\infty = X(\cdot, T) : M^n \rightarrow \mathbb{R}^{n+1}$ with image M_T which coincides with X_∞^l on every ball $B_{R_l}(0)$. Since $M_t \rightarrow M_T$ continuously for $t \rightarrow T$, the smooth convergence holds for $t \rightarrow T$. \square

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

Proof. See [Sch17d, 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. \square

Lemma 4.5 (Hamilton’s trick [Ham86, Lemma 3.5]). *Let $f : [a, b] \times (0, T) \rightarrow \mathbb{R}$ 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,*

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

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

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

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

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

By Lemma 1.4

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

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,

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

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,

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

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

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

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 $|A|_{\max}^2(t_k) \rightarrow \infty$ for $k \rightarrow \infty$. Taking the limit $k \rightarrow \infty$ in (8) yields

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

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

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

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

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

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

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

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

and define

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

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

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

so that

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

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

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

Then there exists a subsequence $((M_\tau^k)_{\tau \in J_k})_{k \in \mathbb{N}}$ that converges on compact subsets of J_∞ 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,

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

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 $x \in \tilde{M}_{\tau_0}^\infty$,

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

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

$$\tilde{M}_{\tau_0}^k \cap B_r^{n+1}(x)$$

can be written as graphs of some function $g_k : B_{r/2}^n(x) \rightarrow \mathbb{R}$ 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,

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

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 $f_k : B_{r/2}^n(x) \times (t_0 - \delta, t_0 + \delta) \rightarrow \mathbb{R}$. Since all time derivatives can be expressed in terms of spatial derivatives, f is smooth in time. By Arzelà–Acsoli, $(f_k)_{k \in \mathbb{N}}$ converges along a subsequence to a smooth limit f_∞ . Like in the proof of Theorem 4.3, we can repeat this process this for a sequence $(R_l)_{l \in \mathbb{N}}$ with $R \geq R_l \rightarrow \infty$, 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). \square

5. TYP-I SINGULARITIES

We want to rescale the surface M_t near a type-I singularity as $t \rightarrow T < \infty$. 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 $k \rightarrow \infty$ 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

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

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

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

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

$$\lambda_k \rightarrow \infty \quad \text{and} \quad \alpha_k \rightarrow -\infty$$

for $k \rightarrow \infty$. Furthermore,

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

where

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

and, for $\delta > 0$,

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

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

Proof.

□

Theorem 5.3 (Convergence of rescalings). *Let $(M_t)_{t \in [0, T]}$ be a smooth, immersed solution of (MCF) with $T < \infty$. Then there exists a sequence of rescaled immersions*

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

that converges for $k \rightarrow \infty$ 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

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

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

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

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 1.7, M_τ^k is embedded for all $k \in \mathbb{N}$ and all $\tau \in [\alpha_k, 0)$. Furthermore,

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

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

□

5.1. Huisken's monotonicity formula.

6. TYP-II SINGULARITIES

The following rescaling technique for type-II singularities was introduced in [Ham95, 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

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

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

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

and define the rescaled embeddings $X_k : M^n \times [\alpha_k, T_k] \rightarrow \mathbb{R}^2$ 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 $X : M^n \times (0, T) \rightarrow \mathbb{R}^2$ be a solution of (MCF) with $T < \infty$. For the type-II rescaling 6.1 in case of a type-II singularity,*

$$\lambda_k \rightarrow \infty, \quad \alpha_k \rightarrow -\infty \quad \text{and} \quad T_k \rightarrow \infty$$

for $k \rightarrow \infty$. Moreover,

$$X_k(0, 0) = 0 \quad \text{and} \quad |A_k|^2(0, 0) = 1$$

for every $k \in \mathbb{N}$ and for any $\varepsilon > 0$ and any $\bar{T} > 0$, 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. □

Theorem 6.3. *Let $(M_t)_{t \in [0, T]}$ be a smooth, immersed solution of (MCF) with $T < \infty$. Then there exists a sequence of rescaled immersions*

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

that converges for $k \rightarrow \infty$ 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

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

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 $\bar{T} \rightarrow \infty$ and $\varepsilon \rightarrow 0$ yields

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

By Proposition 1.7, M_τ^k is embedded for all $k \in \mathbb{N}$ and all $\tau \in [\alpha_k, T_k]$. Furthermore,

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

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

7. 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, φ) , where $U \subset M^n$ is open and $\varphi : U \rightarrow \varphi(U) \subset \mathbb{R}^n$ 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, φ) and (V, ψ) are called C^k -compatible, $k \geq 1$, if

$$\psi \circ \varphi^{-1} : \varphi(U \cap V) \rightarrow \psi(U \cap V)$$

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 $\varphi : U \rightarrow \varphi(U) \subset \mathbb{R}^{n+m}$ 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 ψ and φ .

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 ∂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, φ) of M^n . If M^n is a manifold with boundary, then the interior $\text{int } M^n = M^n \setminus \partial M^n$ is a manifold (without boundary) of dimension n and boundary ∂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 ($\text{rank } DX \equiv n$) with

$$X : M^n \rightarrow \mathbb{R}^{n+m}.$$

We call $M := X(M^n)$ a *hypersurface* in \mathbb{R}^{n+m} . For all $p \in M^n$ and $v, w \in T_p M^n$, the embedding X induces an isomorphism

$$dX_p : T_p M^n \rightarrow T_{X(p)} M,$$

and the *first fundamental form* or *metric* $g_p : T_p M^n \times T_p M^n \rightarrow \mathbb{R}$ 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

$$\partial_k g^{ij} = -g^{pi} g^{qj} \partial_k g_{pq}.$$

The corresponding *Levi-Cevita connection* on M^n is given by

$$\nabla_v w = dX^{-1} \left((D_{dX(v)} dX(w))^\top \right).$$

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 $dX(p)(T_p M^n) = T_{X(p)} M$. The connection can be evaluated in coordinates in terms of the *Christoffel symbols* Γ_{ij}^k defined by

$$\nabla_{\partial_i} \partial_j = \Gamma_{ij}^k \partial_k,$$

where Γ_{ij}^k 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

$$0 = \partial_k \delta_j^i = \partial_k (g^{il} g_{jl}) = g^{il} \partial_k g_{jl} + g_{jl} \partial_k g^{il},$$

so that

$$\begin{aligned} \partial_k g^{ij} &= -g^{il} g^{jm} \partial_k g_{lm} = -g^{il} g^{jm} \partial_k \langle \partial_l X, \partial_m X \rangle \\ &= -g^{il} g^{jm} (\langle \partial_k \partial_l X, \partial_m X \rangle + \langle \partial_l X, \partial_k \partial_m X \rangle) = -g^{il} \Gamma_{kl}^j - g^{jm} \Gamma_{km}^i. \end{aligned}$$

Being in a Levi-Cevita connection the *Li 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 $\operatorname{div}_M : T_p M^n \rightarrow \mathbb{R}$ 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

$$(\operatorname{Hess}_M f)(v, w) := (\nabla^2 f)(v, w) = (\nabla_v f)(w),$$

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* $\Delta_M : C^2(M^n) \rightarrow C^0(M^n)$ is defined as

$$\Delta_M f := \frac{1}{\sqrt{\det g_{kl}}} \partial_j \left(\sqrt{\det g_{kl}} g^{ij} \partial_j f \right) = \operatorname{div}_M (\nabla^M f) = g^{ij} \nabla_i \nabla_j f.$$

We define the *second fundamental form* $\mathbf{A}_p : T_p M^n \times T_p M^n \rightarrow (T_{X(p)} M)^\perp$ by

$$\begin{aligned} \mathbf{A}_p(v, w) &:= - \sum_{k=1}^m \langle D_{dX_p(v)} dX_p(w), \boldsymbol{\nu}_k(p) \rangle \boldsymbol{\nu}_k(p) \\ &= \sum_{k=1}^m \langle dX_p(w), D_{dX_p(v)} \boldsymbol{\nu}_k(p) \rangle \boldsymbol{\nu}_k(p), \end{aligned}$$

where $\{\boldsymbol{\nu}_k\}_{1 \leq k \leq m}$ 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* $\mathbf{H} : M \rightarrow (TM)^\perp$ is the trace of the second fundamental form

$$\mathbf{H} := -g^{ij} \mathbf{A}_{ij} = -g^{ij} \sum_{k=1}^m \langle \partial_i X, \partial_j \nu_k \rangle \nu_k = -\sum_{k=1}^m \operatorname{div}(\nu_k) \nu_k.$$

We calculate that

$$\begin{aligned} \Delta_M X &= g^{ij} (\partial_i \partial_j X - \Gamma_{ij}^k \partial_k X) = g^{ij} \sum_{k=1}^m \langle \partial_i \partial_j X, \nu_k \rangle \nu_k \\ &= -g^{ij} \sum_{k=1}^m \langle \partial_i X, \partial_j \nu_k \rangle \nu_k = \mathbf{H}. \end{aligned}$$

For a submanifold Σ of M , the mean curvature vector is given by

$$\mathbf{H}_\Sigma = -\sum_{k=1}^m \operatorname{div}_\Sigma(\nu_k) \nu_k - \operatorname{div}_\Sigma(\nu_\Sigma) \nu_\Sigma,$$

where ν_Σ is the unit co-normal of Σ . Since ν_Σ tangential to M ,

$$\langle \mathbf{H}_\Sigma, \nu_\Sigma \rangle = -\operatorname{div}_\Sigma \nu_\Sigma$$

and on Σ ,

$$\begin{aligned} \Delta_\Sigma X &= g_\Sigma^{ij} (\partial_i \partial_j X - \Gamma_{ij}^k \partial_k X) \\ &= \sum_{k=1}^m g_\Sigma^{ij} \langle \partial_i \partial_j X, \nu_k \rangle \nu_k + g_\Sigma^{ij} \langle \partial_i \partial_j X, \nu_\Sigma \rangle \nu_\Sigma = \mathbf{H}_\Sigma. \end{aligned}$$

For $m = 1$,

$$\mathbf{A}(v, w) = A(v, w) \nu,$$

where ν is the outward pointing unit normal to M and $A : TM^n \times TM^n \rightarrow \mathbb{R}$ is given by

$$A(v, w) = -\langle D_{dX(v)} dX(w), \nu \rangle = \langle dX(w), D_{dX(v)} \nu \rangle.$$

where ν is the outward pointing unit normal to M . In coordinates,

$$h_{ij} := A(\partial_i, \partial_j) = -\langle \partial_i \partial_j X, \nu \rangle = \langle \partial_i X, \partial_j \nu \rangle.$$

Let $\lambda_1 \leq \lambda_2 \leq \dots \leq \lambda_n$ 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 $k = 1, \dots, n$. The *Weingarten operator* $S : TM^n \rightarrow TM^n$ is given by

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

so that

$$A(v, w) = g(v, S(w)),$$

where in coordinates,

$$h_j^i := g^{ik} h_{kj}$$

and the Weingarten equations by

$$\partial_i \nu = h_i^j \partial_j X. \quad (12)$$

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

$$\mathbf{H} = -g^{ij} h_{ij} \nu = -H \nu,$$

where we define the *mean curvature* H of M as the trace of the second fundamental form with

$$H = g^{ij} h_{ij} = \operatorname{div}_M \nu.$$

We have the Gauss formula

$$\nabla_i \nabla_j X = \partial_i \partial_j X - \Gamma_{ij}^k \partial_k X = -h_{ij} \boldsymbol{\nu} \quad (13)$$

which as before leads to $\Delta_M X = \mathbf{H}$. More useful identities are the Codazzi equations in \mathbb{R}^{n+1}

$$\nabla_k h_{ij} - \nabla_j h_{ik} = \Gamma_{ij}^l h_{lk} - \Gamma_{ik}^l h_{lj} \quad (14)$$

and Simons' identity

$$\Delta h_{ij} = \nabla_i \nabla_j H + H h_{ik} h_j^k - |A|^2 h_{ij}. \quad (15)$$

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

$$R_{kl ij} := g^{kr} R_{lij}^r$$

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

$$R_{ijkl} = h_{ik} h_{jl} - h_{il} h_{jk} \quad (16)$$

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 an open and bounded set such that $M \subset U$. Let $\phi : U \times (-1, 1) \rightarrow U$ 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 $f : U \rightarrow \mathbb{R}^m$ 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 $f : M \rightarrow N$ be differentiable. A point $p \in M$ is called *regular*, if the differential of f in p is surjective. 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 $f : M \rightarrow N$ consists of those points at which the differential $df : TM \rightarrow TN$ has rank less than n as a linear transformation. If $k \geq \max\{n - m + 1, 1\}$, 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 $f : M^m \rightarrow \mathbb{R}^n$ 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 $f : \mathbb{R}^n \rightarrow \mathbb{R}^p$, $f \in C^k$ 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

$$C \subset \bigcup_{i \in \mathbb{N}} W_i \quad \text{and} \quad \sum_{i \in \mathbb{N}} |W_i| < \varepsilon.$$

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 $f : U \rightarrow \mathbb{R}^m$ be Lipschitz. Then $f(C)$ has measure zero.*

Proof. Exercise. □

Definition B.8. A subset C of a differentiable manifold has measure zero, if for every chard $h : U \rightarrow U' \subset \mathbb{R}^m$ 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 , $j = 1, \dots, k$ be numbered so that with $I_j = (a_j, b_j)$ always holds $a_j < a_{j+1}$, $j = 1, \dots, k - 1$. Minimality and covering property imply $a_i < a_{i+1} < b_i < a_{i+2}$. So that

$$\begin{aligned} \sum_i (b_i - a_i) &= \sum_i (a_{i+1} - a_i) + \sum_i (b_i - a_{i+1}) \\ &< \sum_i (a_{i+1} - a_i) + \sum_i (a_{i+1} - a_{i+1}) \leq 2, \end{aligned}$$

where we used that we have telescope sums in the end. □

Theorem B.11 (Fubini). *Let $\mathbb{R}_t^{n-1} := \{x \in \mathbb{R}^n \mid x^n = t\} \subset \mathbb{R}^n$. Let $C \subset \mathbb{R}^n$ be compact and $C_t := C \cap \mathbb{R}_t^{n-1}$ be of measure zero in $\mathbb{R}_t^{n-1} \cong \mathbb{R}^{n-1}$ 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 and attains a positive minimum in the compact set $C \setminus (W_t \times [0, 1])$, which we call α . 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^α 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,

$$\sum_{t_j, i} |W_{t_j}^i \times I_{t_j}^\alpha| \leq 2\varepsilon.$$

Sending $\varepsilon \rightarrow 0$ yields the lemma. \square

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 and 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, $f : U \rightarrow \mathbb{R}^p$ 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

$$D \supset D_1 \supset D_2 \supset \dots$$

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 necessary, 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 $h : U \rightarrow \mathbb{R}^n$ 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 $h : V \rightarrow h(V) = V'$ is a diffeomorphism. Define $g := f \circ h^{-1}$. In a neighbourhood of $h(x)$, g is of the form

$$g : (z^1, \dots, z^n) \mapsto (z^1, g^2(z), \dots, g^n(z)).$$

The hyperplane $\{z \mid z^1 = t\}$ is (locally) mapped into the hyperplane $\{y \mid y^1 = t\}$. Define

$$g_t : \{t\} \times \mathbb{R}^{n-1} \cap V' \rightarrow \{t\} \times \mathbb{R}^{p-1}$$

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 critical 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.

$$\frac{\partial^{k+1} f^1}{\partial x^1 \partial x^{\nu_1} \dots \partial x^{\nu_k}}(x_0) \neq 0.$$

Assume, that this holds in a neighbourhood V of x_0 . Define $w : V \rightarrow \mathbb{R}$ 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

$$h : x \rightarrow (w(x), x^2, \dots, x^n)$$

defines a diffeomorphism $h : V \rightarrow V' = h(V)$. 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

$$g : f \circ h^{-1} : V' \rightarrow \mathbb{R}^p$$

and

$$g_0 : \{0\} \times \mathbb{R}^{n-1} \cap V' \rightarrow \mathbb{R}^p.$$

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

$$|R(x, h)| \leq c|h|^{k+1}$$

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

$$|f(x+h) - f(x)| \leq c \left(\frac{\sqrt{n}a}{r} \right)^{k+1}.$$

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 $r \rightarrow \infty$. \square

Corollary B.13 (Brown). *Let M and N be (finite dimensional) manifolds. Let $f : M \rightarrow N$ be a differentiable (C^∞ -)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 $f : B \rightarrow A$, 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 $f : \overline{B_1(0)} \rightarrow \mathbb{S}^{n-1}$ be a retraction. Show at first, that then there also exists a C^∞ -retraction of $\overline{B_1(0)}$ on \mathbb{S}^{n-1} : We find a retraction g , that is close to $\partial B_1(0)$ of the class C^∞ , e.g.,

$$g(x) = \begin{cases} f\left(\frac{x}{|x|}\right) & \text{for } \frac{1}{2} \leq |x| \leq 1 \\ f(2x) & \text{for } 0 \leq |x| \leq \frac{1}{2}. \end{cases}$$

Mollification in the interior gives a C^∞ -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 $y, z \in f^{-1}(y)$. \square

Theorem B.16 (Brouwer's fixed point theorem). *Let $f : \overline{B_1(0)} \rightarrow \overline{B_1(0)}$ 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} . \square

APPENDIX C. MAXIMUM PRINCIPLES

We will need the following maximum principles.

Theorem C.1 (Strong elliptic maximum principle). *Let M be closed and $f : M \rightarrow \mathbb{R}$ satisfy*

$$-\Delta_M f + b^i \nabla_i^M f + cf \leq 0$$

for some smooth funtions b^i and $c \leq 0$. If $f \leq 0$, but $f \not\equiv 0$, then $f < 0$.

Proof. For a proof see [Eva02, §6.4, Theorem 4] or [Sch17b, Theorem 5.5] for $M^n = \mathbb{R}^n$. \square

Let M^n be a smooth n -dimensional manifold with boundary whose closure is compact. Let $X : \bar{M}^n \times [0, T] \rightarrow \mathbb{R}^{n+m}$ 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

$$L(f) := \partial_t f - a^{ij} \nabla_i \nabla_j f - b^i \nabla_i f - cf, \quad (17)$$

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,

$$\lambda |\xi|^2 \leq a_{ij} \xi_i \xi_j \leq \Lambda |\xi|^2 \quad (18)$$

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

$$\begin{aligned} U_R(p_0, t_0) &:= X^{-1}(B_R(X(p_0, t_0)) \cap M_{t_0}) \\ &= \{p \in M^n \mid |X(p, t_0) - X(p_0, t_0)| < R\}, \end{aligned}$$

the parabolic neighbourhood

$$\begin{aligned} Q_R(p_0, t_0) &:= \{(p, t) \in M^n \times (t_0 - R^2, t_0] \mid |X(p, t) - X(p_0, t)| < R\} \\ &= \bigcup_{t \in (t_0 - R^2, t_0]} (U_R(p_0, t) \times \{t\}) \end{aligned}$$

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 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 $L(f) \geq 0$ on Q and $f \geq 0$ on $\mathcal{P}Q$. Then $f \geq 0$ in Q .*
- (ii) *If $L(f) \leq 0$ on Q and $f \leq 0$ on $\mathcal{P}Q$. Then $f \leq 0$ in Q .*

Proof.

□

To prove a strong maximum principle, we need the following theorem.

Theorem C.3 (Weak Harnack inequality). *Let $U_R(p_0, t_0) \times [t_1, t_0] \subset M^n \times [0, T]$ and $f \in C^{2;1}(U_R(p_0, t_0) \times [t_1, t_0]) \cap C^0(\bar{U}_R(p_0, t_0) \times [t_1, t_0])$ be a non-negative solution of*

$$L(f) \geq 0$$

in $U_R(p_0, t_0) \times [t_1, t_0]$ where there exists $\lambda, \Lambda, \nu > 0$ such that

$$\lambda g^{ij} \leq a^{ij} \leq \Lambda g^{ij} \quad \text{and} \quad \left(\frac{|b|^2}{\lambda} + |c| \right) \leq \frac{\nu \lambda}{R^2}$$

for all $t \in [t_1, t_0]$. Let

$$f \geq h > 0$$

in $U_{\varepsilon R}(p_0, t_0) \times \{t_1\}$ for some $0 < \varepsilon < 1/2$. Then

$$f \geq \varepsilon^\kappa \frac{h}{2}$$

in $U_{R/2}(p_0, t_0) \times \{t_0\}$ for $\kappa = \kappa(n, R, t_0 - t_1, \lambda, \Lambda, \nu, X_t, \frac{\partial}{\partial p_i} X_t, \frac{\partial^2}{\partial p_i \partial p_j} X_t) > 0$, where X_t is the embedding of M_t .

Proof.

□

Theorem C.4 (Strong 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 $L(f) \geq 0$ 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 $(p_0, t_0) \in Q \setminus \mathcal{P}Q$ with $f(p_0, t_0) = \max_{\bar{Q}} f$, then f is constant in \bar{Q} .*

Proof. □

Theorem C.5 (Strong maximum principle for parabolic equations). *Let M be closed and $f : M \times [0, T) \rightarrow \mathbb{R}$ satisfy*

$$\partial_t f \geq \Delta_{M_t} f + b^i \nabla_i^{M_t} f + c f$$

for some smooth functions b^i and $c \geq 0$. If $f(\cdot, 0) \geq 0$ then $\min_M f(\cdot, t) \geq \min_M f(\cdot, 0)$. Furthermore, if $f(p, t_0) = \min_M f(\cdot, 0)$ for some $p \in M$, $t > 0$, then $f = \min_M f(\cdot, 0)$ for $0 \leq t \leq t_0$.

Proof. See for example [Eva02, Chapters 6.4 and 7.1.4] or [Sch17c, Theorem 3.7] for $M^n = \mathbb{R}^n$. □

Let m_{ij} be a symmetric tensor. We say that m_{ij} is nonnegative, $m_{ij} \succeq 0$, if all eigenvalues of m_{ij} are nonnegative. Let u^k be a vector field and let g_{ij} , m_{ij} and n_{ij} be symmetric tensors on a compact manifold M which may all depend on time. Assume that $n_{ij} = p(m_{ij}, g_{ij})$ is a polynomial in m_{ij} formed by contracting products of m_{ij} with itself using the metric. Furthermore, let this polynomial satisfy a null-eigenvector condition, i.e. for any null-eigenvector v of m_{ij} we have $n_{ij} v^i v^j \geq 0$. Then we have

Theorem C.6 (Weak parabolic maximum principle for symmetric 2-tensors I, Hamilton, [Ham82, Theorem 9.1]). *Suppose that on $[0, T)$ the evolution equation*

$$\partial_t m_{ij} = \Delta_{M_t} m_{ij} + u^k \nabla_k^{M_t} m_{ij} + n_{ij}$$

holds, where $n_{ij} = p(m_{ij}, g_{ij})$ satisfies the null-eigenvector condition above. If $m_{ij} \succeq 0$ at $t = 0$, then it remains so on $(0, T)$.

Proof. □

Exercise C.7. Proof Theorems C.1, C.5 and C.6 on Riemannian manifolds.

There is also a strong version of the maximum principle for tensors. We need the following lemma.

Lemma C.8 (Hamilton, [Ham86, Lemma 8.1]). *Let V a vector bundle over a compact manifold M^n . Let v be a smooth section of V satisfying*

$$\frac{\partial v}{\partial t} = \Delta_{M_t} v + \phi(v).$$

Let $f(v)$ be a convex function on the bundle invariant under parallel translation whose level curves $f(v) \leq c$ are preserved by the ODE

$$\frac{dv}{dt} = \phi(v).$$

Then the inequality $f(v) \leq c$ is preserved by the PDE for any constant c . Furthermore if $f(v) < c$ at one point at time $t = 0$, then $f(v) < c$ everywhere on M for all $t > 0$.

Proof. □

Theorem C.9 (Strong parabolic maximum principle for symmetric 2-tensors II, Hamilton, [Ham86, Lemma 8.2]). *Let M^n be closed and m_{ij} be a symmetric bilinear form of a vector bundle V over M^n , which solves*

$$\partial_t m_{ij} = \Delta_{M_t} m_{ij} + \phi_{ij},$$

where ϕ_{ij} is a symmetric bilinear form, depending on m_{ij} , with the property $\phi_{ij} \succeq 0$ if $m_{ij} \succeq 0$. If $m_{ij} \succeq 0$ for $t = 0$ then $m_{ij} \succeq 0$ for all $t \leq 0$. Furthermore, for $t > 0$, the rank of the null-space of m_{ij} is constant, and the null-space is invariant under parallel transport and invariant in time.

Proof. □

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