

On Conformally Flat Submanifolds of Codimension 2 with Parallel Mean Curvature Vector in a Euclidean Space

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0. Introduction. A Riemannian manifold (M, g) is said to be conformally flat, if it is locally conformally equivalent to a Euclidean space.

The study of conformally flat hypersurfaces in a conformally flat space arose from Schouten [12], who showed that an $n(\geq 4)$ -dimensional hypersurface M in a Euclidean space is conformally flat if and only if M is quasi-umbilic. Afterward this subject has been investigated by many authors, for example, [3], [4], and [10]. However, Cecil and Ryan [1] showed the counter example against [10].

Conformally flat submanifolds of codimension 2 is treated by Chen and Yano [5]. On the other hand, Moore [8] and Moore and Morvan [9] studied conformally flat submanifolds whose codimension is greater than 2. Especially, the latter studied the second fundamental form of such submanifold. And Kitagawa [7] studied the shape of such submanifolds.

The purpose of this paper is to study locally the shape of a conformally flat manifold M^n ($n \geq 5$) which is isometrically immersed into a Euclidean space with parallel mean curvature vector. Then we obtain the next theorem:

THEOREM. *Let M^n ($n \geq 5$) be a conformally flat manifold which is isometrically immersed into an $(n+2)$ -dimensional Euclidean space E^{n+2} . Assume that the mean curvature vector is parallel and it is non-trivial on M . Then if we take off points of constant sectional curvature and a certain subset of measure 0 from M , the rest satisfies locally one of the following two conditions, where (x^1, \dots, x^{n+2}) denotes the natural coordinate of E^{n+2} :*

(1) *There exist a curve $(x^1(t), x^2(t), 0, \dots, 0)$ in E^{n+2} and a C^∞ -function $\gamma(t)$ (> 0) such that M is denoted by*

$$(x^1(t), x^2(t), x^3, \dots, x^{n+2}),$$

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where

$$(x^3)^2 + \cdots + (x^{n+2})^2 = \gamma(t)^2.$$

(2) There exist a surface $(x^1(u, v), x^2(u, v), x^3(u, v), 0, \cdots, 0)$ in E^{n+2} and a C^∞ -function $\gamma(u, v) (>0)$ such that M is denoted by

$$(x^1(u, v), x^2(u, v), x^3(u, v), x^4, \cdots, x^{n+2}),$$

where

$$(x^4)^2 + \cdots + (x^{n+2})^2 = \gamma(u, v)^2.$$

In Section 1, the notations are explained and some properties of the manifold M which we will seek are prepared for. Section 2 is devoted to the study of the situation of the second fundamental form of M , which is divided into different three cases (Theorem 2.1). In particular, M is totally umbilic in a Euclidean space in the case of III in Theorem 2.1. In Section 3, we study the shape of M in the case of II in Theorem 2.1 and obtain that it becomes as the statement (1). In Section 4, we study the shape of M in the case of I in Theorem 2.1 and obtain that it becomes as the statement (2).

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1. Preliminaries. Let M be an $n(\geq 2)$ -dimensional connected Riemannian manifold with the Riemannian metric g . Let f be an isometric immersion of M into an $(n+p)$ -dimensional Euclidean space E^{n+p} . Since the argument is only considered in the local version, M needs not to be distinguished from $f(M)$. So, in order to simplify the discussion, we identify a point q in M with the point $f(q)$ and a tangent vector X at q with the tangent vector $df_q X$.

Now, we choose an orthonormal local frame field $\{e_1, \cdots, e_n, e_{n+1}, \cdots, e_{n+p}\}$ on E^{n+p} in such a way that, restricted to M , the vectors e_1, \cdots, e_n are tangent to M , and hence the others are normal to M . With respect to the fields of frame of E^{n+p} , let $\{\omega^1, \cdots, \omega^{n+p}\}$ be the dual forms. Here and in the sequel, the following convention on the range of the indices is adopted, unless otherwise stated:

$$A, B, \cdots = 1, \cdots, n, n+1, \cdots, n+p,$$

$$i, j, \cdots = 1, \cdots, n,$$

$$\alpha, \beta, \cdots = n+1, \cdots, n+p.$$

Then, associated to the frame field $\{e_A\}$, there exist 1-forms $\tilde{\omega}_B^A$ on E^{n+p} so that they satisfy the following structure equations on E^{n+p} :

$$\begin{aligned}
 (1.1) \quad & d\tilde{\omega}^A + \tilde{\omega}_B^A \wedge \tilde{\omega}^B = 0, \\
 & \tilde{\omega}_B^A + \tilde{\omega}_A^B = 0, \\
 & d\tilde{\omega}_B^A + \tilde{\omega}_C^A \wedge \tilde{\omega}_B^C = 0,
 \end{aligned}$$

where the Einstein convention for the summation is adopted. Let D denote the connection on E^{n+p} . Then we have

$$D_X e_A = \tilde{\omega}_A^B(X) e_B$$

for any tangent vector X on E^{n+p} . And $\{\tilde{\omega}_B^A\}$ are called *connection forms* on E^{n+p} .

Restricting $\{\tilde{\omega}^A\}$ and $\{\tilde{\omega}_B^A\}$ on E^{n+p} to the submanifold M , we denote them by ω^A and ω_B^A respectively, that is,

$$\begin{aligned}
 \tilde{\omega}^A|_M &= \omega^A, \\
 \tilde{\omega}_B^A|_M &= \omega_B^A.
 \end{aligned}$$

It then yields

$$(1.2) \quad \omega^\alpha = 0.$$

Let ∇ denote the connection of M . Then $\nabla_X Y$ is equal to the tangential component of $D_X Y$ to M , where X and Y are any tangent vector fields on M . The metric g in M induced from the standard metric in the ambient space E^{n+p} is given by

$$g = \sum_{i=1}^n \omega^i \otimes \omega^i$$

in other words,

$$g = \sum_{i=1}^n (\omega^i)^2.$$

Then $\{e_1, \dots, e_n\}$ are also the orthonormal frame field with respect to g , and $\{\omega^1, \dots, \omega^n\}$ are the dual fields with respect to $\{e_1, \dots, e_n\}$. It follows from (1.1) and the Cartan's lemma that

$$\begin{aligned}
 (1.3) \quad & \omega_i^\alpha = h_{ij}^\alpha \omega^j, \\
 & h_{ij}^\alpha = h_{ji}^\alpha.
 \end{aligned}$$

The quadratic form $h_{ij}^\alpha \omega^i \otimes \omega^j$ is called the *second fundamental form* of the immersion f on M in the direction of e_α . The second fundamental form α of M can be written as

$$(1.4) \quad \alpha(X, Y) = h_{ij}^\beta \omega^i(X) \omega^j(Y) e_\beta$$

for any tangent vectors X and Y on M .

From the structure equations (1.1) of the ambient space the following structure equations on the submanifold M are given :

$$(1.5) \quad \begin{aligned} d\omega^i + \omega_j^i \wedge \omega^j &= 0, \\ d\omega_j^i + \omega_k^i \wedge \omega_j^k &= \Omega_j^i, \\ \Omega_j^i &= -\frac{1}{2} R_{jkh}^i \omega^k \wedge \omega^h, \end{aligned}$$

where ω_j^i and Ω_j^i denote the connection form and the curvature form on the submanifold M , respectively. Moreover they yield

$$(1.6) \quad \begin{aligned} d\omega_\beta^\alpha + \omega_\gamma^\alpha \wedge \omega_\beta^\gamma &= \Omega_\beta^\alpha \\ \Omega_\beta^\alpha &= -\frac{1}{2} R^{N\alpha}{}_{\beta kh} \omega^k \wedge \omega^h, \end{aligned}$$

where ω_β^α defines the connection form induced in the normal bundle $N(M)$ of M and Ω_β^α is called the *normal curvature form* of M and R^N is called the *normal curvature*. If $\Omega_\beta^\alpha = 0$ for any α, β , the normal connection is said to be *flat*.

Taking the exterior derivative of (1.3), we have

$$(1.7) \quad (dh_{ij}^\alpha - h_{ik}^\alpha \omega_j^k - h_{kj}^\alpha \omega_i^k + h_{ij}^\beta \omega_\beta^\alpha) \wedge \omega^j = 0.$$

So, we define the covariant derivative h_{ij}^α of h_{ij}^α by

$$(1.8) \quad h_{ij}^\alpha \omega^k = dh_{ij}^\alpha - h_{ik}^\alpha \omega_j^k - h_{kj}^\alpha \omega_i^k + h_{ij}^\beta \omega_\beta^\alpha.$$

It turns out that (1.7) says that

$$h_{ij}^\alpha \omega^j \wedge \omega^k = 0,$$

which is equivalent to

$$h_{ijk}^\alpha = h_{ikj}^\alpha.$$

The mean curvature vector η on the submanifold M is defined by

$$\eta = \frac{1}{n} \sum_i \alpha(e_i, e_i) = \frac{1}{n} \sum_i h_{ii}^\beta e_\beta,$$

which is independent of the choice of the orthonormal frame field $\{e_i\}$. The mean curvature vector η is said to be *parallel* if $D_X \eta$ is tangent to M for any tangent vector X on M . It is equivalent that

$$\sum_i h_{ii}^\alpha = 0$$

for all indices k and α .

For each normal vector ξ at x , a linear transformation A_ξ on the tangent space $T_x(M)$ is defined by

$$g(A_\xi X, Y) = g(\alpha(X, Y), \xi)$$

for any X and $Y \in T_x(M)$. A_ξ is symmetric with respect to the metric g and it is called the *shape operator* or the *second fundamental tensor* with respect to the normal vector ξ . In particular if $\xi = e_\beta$, then

$$g(A_\beta X, Y) = h_{ij}^{\beta} \omega^i(X) \omega^j(Y),$$

where $A_\beta = A_{e_\beta}$.

Now, a Riemannian manifold is said to be *conformally flat* if each point of M has a neighborhood where there exists a conformal diffeomorphism onto a subset in a Euclidean space.

For a conformally flat submanifold immersed into a Euclidean space, Moore and Morvan [9] proved the following property:

LEMMA 1.1. *Let M be an n -dimensional conformally flat submanifold in E^{n+p} . If $p \leq 4$ and $p \leq n-3$, then at each point x of M there exists a normal vector ξ such that*

$$\langle \beta(X, Y), \beta(Z, W) \rangle = \langle \beta(X, W), \beta(Z, Y) \rangle$$

for any vectors X, Y, Z and W at x , where

$$\beta(X, Y) = \alpha(X, Y) - \langle X, Y \rangle \xi$$

and α is the second fundamental form of the submanifold and $\langle \cdot, \cdot \rangle$ is the metric of E^{n+p} , and hence of M .

DEFINITION 1.2. M is said to be *quasi-umbilic* in the sense of Moore and Morvan if at each point of M there exist orthonormal normal vectors e_α ($\alpha = n+1, \dots, n+p$) such that the second fundamental tensor A_α with respect to each normal vector e_α has only two distinct eigenvalues with multiplicity $n-1$ and 1 or n and 0.

THEOREM 1.3. (Moore and Morvan [9]) *Let M be an n -dimensional submanifold in E^{n+p} . If $p \leq 4$ and $p \leq n-3$ and if M is conformally flat, then M is quasi-umbilic in the sense of Moore and Morvan.*

We must remark that in Lemma 1.1. it is not proved for the normal vector ξ to be smooth and in Theorem 1.3. it is also not proved for e_α ($\alpha = n+1, \dots, n+p$) to be smooth. For the smoothness of ξ Kitagawa [7] asserted the followings:

LEMMA 1.4. (Kitagawa [7]) *If the same assumption as Lemma 1.1. is satisfied, then there exists an open dense set M^* of M on which ξ is smooth.*

THEOREM 1.5. (Kitagawa [7]) Suppose \mathcal{D} be a distribution on M^* defined by

$$\mathcal{D}(q) = \{X \in T_q(M) : \alpha(X, Y) - \langle X, Y \rangle \xi = 0 \text{ for all } Y \in T_q(M)\}.$$

Then \mathcal{D} is completely integrable and its integral manifold is umbilic in E^{n+p} .

From now on, we study M^* in place of M . Lemma 1.4. implies that the complement of M^* in M is of measure 0.

2. Conformally flat submanifolds of codimension 2 with parallel mean curvature vector in a Euclidean space.

From now on, we assume M^n is isometrically immersed into $(n+2)$ -dimensional Euclidean space E^{n+2} . Let η be the mean curvature vector on the submanifold.

THEOREM 2.1. Let M be a conformally flat submanifold in E^{n+2} . If the mean curvature vector η is parallel and non-trivial on M , then there exists an open dense subset M^* in M as follows:

At each point p of M^* , there are a neighborhood U of p , a C^∞ orthonormal frame field e_1, \dots, e_n and C^∞ orthonormal normal vector fields e_{n+1}, e_{n+2} such that the connection form with respect to $e_1, \dots, e_n, e_{n+1}, e_{n+2}$ satisfies one of the following three conditions I, II, III on U :

$$\begin{aligned} \text{I} & \begin{cases} \omega_s^{n+1} = \lambda \omega^s, & \omega_s^{n+2} = \tau \omega^s, \\ \omega_{n-1}^{n+1} = \mu \omega^{n-1}, & \omega_{n-1}^{n+2} = \tau \omega^{n-1}, \\ \omega_n^{n+1} = \lambda \omega^n, & \omega_n^{n+2} = \rho \omega^n, \\ \text{where } \lambda \neq \mu, \tau \neq \rho \text{ and } s=1, \dots, n-2. \end{cases} \\ \text{II} & \begin{cases} \omega_a^{n+1} = \lambda \omega^a, & \omega_a^{n+2} = \tau \omega^a, \\ \omega_n^{n+1} = \lambda \omega^n, & \omega_n^{n+2} = \rho \omega^n, \\ \text{where } \tau \neq \rho \text{ and } a=1, \dots, n-1. \end{cases} \\ \text{III} & \begin{cases} \omega_i^{n+1} = \lambda \omega^i, & \omega_i^{n+2} = \tau \omega^i, \\ \text{where } i=1, \dots, n. \end{cases} \end{aligned}$$

PROOF. Choose arbitrary C^∞ orthonormal normal vector fields $\tilde{e}_{n+1}, \tilde{e}_{n+2}$. Since the codimension is 2 and the mean curvature vector η is parallel and non-trivial, which implies $\eta/\|\eta\|$ is parallel, it turned out that the unit normal vector which is orthonormal to η must be also parallel.

This yields that the normal connection is flat. It follows from the flatness of the normal connection that the second fundamental tensors of \tilde{e}_{n+1} and \tilde{e}_{n+2} are simultaneously diagonalizable. Therefore, at each point p of M there exist a neighborhood V of p and an orthonormal C^∞ frame e_1, \dots, e_n on V such that the second fundamental tensors $\tilde{A}_{n+1}, \tilde{A}_{n+2}$ with respect to $\tilde{e}_{n+1}, \tilde{e}_{n+2}$ are represented as follows:

$$\langle \tilde{A}_{n+1}(e_i), e_j \rangle = \begin{pmatrix} \nu_1 & & & \\ & \nu_2 & & \\ & & \ddots & \\ & & & \nu_n \end{pmatrix}$$

$$\langle \tilde{A}_{n+2}(e_i), e_j \rangle = \begin{pmatrix} \pi_1 & & & \\ & \pi_2 & & \\ & & \ddots & \\ & & & \pi_n \end{pmatrix}$$

where ν_i, π_i ($i=1, \dots, n$) are C^∞ -functions on V .

Let α be the second fundamental form. Then,

$$\begin{aligned} \alpha(e_i, e_j) &= \langle \tilde{A}_{n+1}(e_i), e_j \rangle e_{n+1} + \langle \tilde{A}_{n+2}(e_i), e_j \rangle e_{n+2} \\ &= \nu_i \delta_{ij} e_{n+1} + \pi_i \delta_{ij} e_{n+2}. \end{aligned}$$

Therefore

$$\beta(e_i, e_j) = \nu_i \delta_{ij} e_{n+1} + \pi_i \delta_{ij} e_{n+2} - \delta_{ij} \xi,$$

where β and ξ are those defined in Lemma 1.1. Clearly,

$$(2.1) \quad \beta(e_i, e_j) = 0 \quad \text{if } i \neq j.$$

Replacing the indices if necessary, we may consider

$$(2.2) \quad \begin{aligned} \beta(e_l, e_l) &= 0, \dots, \beta(e_l, e_l) = 0, \\ \beta(e_{l+1}, e_{l+1}) &\neq 0, \dots, \beta(e_n, e_n) \neq 0 \end{aligned}$$

at each point of V . We set $r = n - l$, and consider r as a function on V . If $r(q) > 0$ for a point $q \in V$, then there exists a neighborhood $U (\subset V)$ of q where r is equal to $r(q)$.

For each t ($t=l+1, \dots, n$) we define a real-valued function c_t and a normal vector field ζ_t by

$$(2.3) \quad \begin{aligned} \beta(e_t, e_t) &= c_t \zeta_t, \\ \langle \zeta_t, \zeta_t \rangle &= 1. \end{aligned}$$

From (2.2), we find

$$(2.4) \quad c_t \neq 0 \quad \text{for } t=l+1, \dots, n:$$

Making use of (2.3), (2.4) and Lemma 1.1., we get

$$\langle \zeta_u, \zeta_t \rangle = \frac{1}{c_u c_t} \langle \beta(e_u, e_t), \beta(e_u, e_t) \rangle$$

for $u, t = l+1, \dots, n$. Therefore

$$(2.5) \quad \langle \zeta_u, \zeta_t \rangle = 0 \quad \text{if } u \neq t$$

because of (2.1). By (2.3) and (2.5), we find ζ_t ($t = l+1, \dots, n$) are orthonormal normal vector fields.

From definition,

$$\zeta_t = \frac{\beta(e_t, e_t)}{\|\beta(e_t, e_t)\|},$$

where $\|\ \ \|$ denotes the norm with respect to the metric $\langle \ , \ \rangle$. Since e_t is a C^∞ unit vector field on $U \subset V$, ζ_t is a C^∞ vector field on U . By the assumption that the codimension is 2, r ($= n-l$) is at most 2.

First we consider the case when $r=2$ on U . Using (2.1), (2.2) and (2.3), we get

$$(\beta(e_i, e_j))_{i,j=1,\dots,n} = \begin{pmatrix} 0 & & & & & \\ & 0 & & & & \\ & & \dots & & & \\ & & & 0 & & \\ & & & & c_{n-1}\zeta_{n-1} & \\ & & & & & c_n\zeta_n \end{pmatrix}$$

where $c_{n-1} \neq 0$ and $c_n \neq 0$. By the definition of β in Lemma 1.1., we find

$$\alpha(e_i, e_j) = \beta(e_i, e_j) + \langle e_i, e_j \rangle \xi.$$

Considering ζ_{n-1}, ζ_n as e_{n+1}, e_{n+2} , we have

$$\begin{aligned} \omega_s^{n+1} &= \langle \xi, \zeta_{n-1} \rangle \omega^s, \\ \omega_{n-1}^{n+1} &= (c_{n-1} + \langle \xi, \zeta_{n-1} \rangle) \omega^{n-1}, \\ \omega_n^{n+1} &= \langle \xi, \zeta_{n-1} \rangle \omega^n, \\ \omega_s^{n+2} &= \langle \xi, \zeta_n \rangle \omega^s, \\ \omega_{n-1}^{n+1} &= \langle \xi, \zeta_n \rangle \omega^{n-1}, \\ \omega_n^{n+2} &= (c_n + \langle \xi, \zeta_n \rangle) \omega^n, \end{aligned}$$

where $s=1, \dots, n-2$. Therefore, we find the condition I is satisfied on U .

Next, we study the case where $r=1$ on U . Since ζ_n only exists, we define the unit normal vector ζ_{n-1} such that ζ_{n-1} is orthogonal to ζ_n . By (2.1), (2.2) and (2.3), we get

$$(\beta(e_i, e_j))_{i,j=1,\dots,n} = \begin{pmatrix} 0 & & & \\ & \ddots & & \\ & & \ddots & \\ & & & 0 \\ & & & & c_n \zeta_n \end{pmatrix},$$

where $c_n \neq 0$. By definition of β , we find

$$\alpha(e_i, e_j) = \beta(e_i, e_j) + \langle e_i, e_j \rangle \xi.$$

Considering ζ_{n-1}, ζ_n as e_{n+1}, e_{n+2} respectively, we have

$$\begin{aligned} \omega_a^{n+1} &= \langle \xi, \zeta_{n-1} \rangle \omega^a, \\ \omega_n^{n+1} &= \langle \xi, \zeta_{n-1} \rangle \omega^n, \\ \omega_a^{n+2} &= \langle \xi, \zeta_n \rangle \omega^a, \\ \omega_n^{n+2} &= (c_n + \langle \xi, \zeta_n \rangle) \omega^n, \end{aligned}$$

where $a=1, \dots, n-1$. Therefore we find the condition II is satisfied on U .

Now, we define a subset \tilde{V} in V by

$$\tilde{V} = \{q \in V : r(q) = 0\}^o,$$

where A^o denote the set of all inner points of a set A . On \tilde{V}

$$\beta(e_i, e_j) = 0.$$

Let e_{n+1}, e_{n+2} be arbitrary C^∞ orthonormal normal vector fields. Then we get

$$\begin{aligned} (\langle \alpha(e_i, e_j), e_{n+1} \rangle)_{i,j=1,\dots,n} &= \begin{pmatrix} \langle \xi, e_{n+1} \rangle & & & \\ & \ddots & & \\ & & \ddots & \\ & & & \langle \xi, e_{n+1} \rangle \end{pmatrix} \\ (\langle \alpha(e_i, e_j), e_{n+2} \rangle) &= \begin{pmatrix} \langle \xi, e_{n+2} \rangle & & & \\ & \ddots & & \\ & & \ddots & \\ & & & \langle \xi, e_{n+2} \rangle \end{pmatrix} \end{aligned}$$

on \tilde{V} . From which it follows that the condition III is satisfied, where we consider \tilde{V} as U . Q. E. D.

REMARK 2.2. In the case of III, U is totally umbilic in E^{n+2} . It is then known that this implies that there exists a totally geodesic $(n+1)$ -dimensional plane L in E^{n+2} such that U is a part of n -dimensional sphere of L .

REMARK. 2.3. We call $p \in M$ a *point of constant sectional curvature* if all sectional curvatures at p have the same value. When we study the shape of U which satisfies I or II, we avoid the points of constant sectional curvature.

3. The case of II in Theorem 2.1.

In this section we study the shape of U on which the condition II in Theorem 2.1. is satisfied.

In this case there exist C^∞ orthonormal normal vector fields e_{n+1}, e_{n+2} and a C^∞ orthonormal frame e_1, \dots, e_n on U such that

$$(3.1) \quad \begin{aligned} \omega_a^{n+1} &= \lambda \omega^a, & \omega_a^{n+2} &= \tau \omega^a, \\ \omega_n^{n+1} &= \lambda \omega^n, & \omega_n^{n+2} &= \rho \omega^n, \end{aligned}$$

where λ, τ and ρ are C^∞ -functions on U , $\tau \neq \rho$ on U and $a=1, \dots, n-1$.

In this section indices A, B, \dots run over the range $\{1, \dots, n+2\}$ and a, b, \dots move from 1 to $n-1$. For indices a, b, \dots the Einstein's convention is used.

LEMMA 3.1. *A point p in M is of constant sectional curvature if and only if $\tau=0$ at p .*

PROOF. Let $K(X, Y)$ denote the sectional curvature of the plane spanned by $X, Y \in T_p(M)$. Then

$$\begin{aligned} K(e_a, e_b) &= \lambda^2 + \tau^2, \\ K(e_a, e_n) &= \lambda^2 + \tau\rho, \end{aligned}$$

which imply that if p is a point of constant sectional curvature, then $\tau=0$.

Conversely let $\tau=0$. Let X, Y be any unit vectors which are mutually orthogonal. If we set

$$\begin{aligned} X &= \sum_{i=1}^n X^i e_i, \\ Y &= \sum_{j=1}^n Y^j e_j, \end{aligned}$$

then

$$\begin{aligned} \sum_{i=1}^n (X^i)^2 &= 1, \\ \sum_{j=1}^n (Y^j)^2 &= 1, \end{aligned}$$

$$\sum_{i=1}^n X^i Y^i = 0,$$

Therefore

$$K(X, Y) = \lambda^2. \quad \text{Q. E. D.}$$

Since we have assumed that U has no point of constant sectional curvature, we may consider $\tau \neq 0$ on U .

We set

$$\begin{aligned} d\lambda &= \lambda_a \omega^a + \lambda_n \omega^n, \\ (3.2) \quad d\tau &= \tau_a \omega^a + \tau_n \omega^n, \\ d\rho &= \rho_a \omega^a + \rho_n \omega^n. \end{aligned}$$

LEMMA 3.2.

$$\begin{aligned} (3.3) \quad \omega_n^a &= -\frac{\tau_n}{\tau - \rho} \omega^a, \\ \omega_{n+2}^{n+1} &= 0, \end{aligned}$$

$$(3.4) \quad \lambda_a = 0, \quad \tau_a = 0, \quad \rho_a = 0, \quad \lambda_n = 0.$$

PROOF. We set

$$(3.5) \quad \omega_n^a = A_b^a \omega^b + A_n^a \omega^n,$$

$$(3.6) \quad \omega_{n+2}^{n+1} = C_b \omega^b + C_n \omega^n.$$

Taking the exterior derivative of the first equation of (3.1) and using (3.2), (3.5), (3.1) and (3.6), we get

$$\begin{aligned} d(\lambda \omega^a) &= (\lambda_b \omega^b + \lambda_n \omega^n) \wedge \omega^a + \lambda \omega^b \wedge \omega_b^a + \lambda \omega^n \wedge A_b^a \omega^b, \\ d\omega_a^{n+1} &= \lambda \omega^b \wedge \omega_b^a - \lambda A_b^a \omega^b \wedge \omega^n + \tau C_b \omega^a \wedge \omega^b + \tau C_n \omega^a \wedge \omega^n. \end{aligned}$$

This implies that

$$\lambda_b + \tau C_b = 0,$$

$$\lambda_n + \tau C_n = 0.$$

Similarly taking the exterior derivative of the other equations of (3.1), we obtain

$$\lambda_b + \rho C_b = 0,$$

$$\tau_b - \lambda C_b = 0,$$

$$\tau_n \delta_b^a + (\tau - \rho) A_b^a - \lambda C_n \delta_b^a = 0,$$

$$-\rho_b + (\rho - \tau) A_n^b + \rho C_b = 0.$$

Since $\tau \neq \rho$ and $\tau \neq 0$, we obtain

$$C_b=0, \quad \lambda_b=0, \quad \tau_b=0, \quad C_n=-\frac{\lambda_n}{\tau},$$

$$A_n^b=-\frac{\rho_b}{\tau-\rho}, \quad A_b^a=\frac{\lambda C_n-\tau_n}{\tau-\rho}\delta_b^a.$$

Therefore

$$\omega_n^a=-\frac{\lambda\lambda_n+\tau\tau_n}{\tau(\tau-\rho)}\omega^a-\frac{\rho_a}{\tau-\rho}\omega^n,$$

$$\omega_{n+2}^{n+1}=-\frac{\lambda_n}{\tau}\omega^n,$$

$$\lambda_b=0, \quad \tau_b=0.$$

On the other hand, by (3.1), we get

$$\eta=\lambda e_{n+1}+\frac{1}{n}\{(n-1)\tau+\rho\}e_{n+2},$$

where η is the mean curvature vector. The condition that the mean curvature vector η is parallel implies

$$\lambda_a+\frac{1}{n}\{(n-1)\tau+\rho\}C_a=0,$$

$$-\lambda_a C+\frac{1}{n}\{(n-1)\tau_a+\rho_a\}=0,$$

$$\lambda_n+\frac{1}{n}\{(n-1)\tau+\rho\}C_n=0,$$

$$-\lambda C_n+\frac{1}{n}\{(n-1)\tau_n+\rho_n\}=0.$$

Using the fact that

$$\lambda_a=0, \quad C_a=0, \quad \tau_a=0,$$

$$\lambda_n+\tau C_n=0,$$

we obtain

$$\rho_a=0, \quad C_n=0, \quad \lambda_n=0.$$

Therefore

$$\omega_n^a=-\frac{\tau}{\tau-\rho}\omega^a,$$

$$\omega_{n+2}^{n+1}=0,$$

$$\lambda_a=0, \quad \tau_a=0, \quad \rho_n=-(n-1)\tau_n.$$

Q. E. D.

We set

$$d\tau_n=\tau_{na}\omega^a+\tau_{nn}\omega^n,$$

$$d\rho_n = \rho_{na}\omega^a + \rho_{nn}\omega^n.$$

Because of $d^2=0$, we obtain

$$(3.7) \quad \tau_{na} = 0, \quad \rho_{na} = 0.$$

By (3.3), we get

$$-(\tau - \rho)\omega_n^a = \tau_n\omega^a.$$

By taking the exterior derivative of the above, we get

$$(3.8) \quad \tau_{nn}(\tau - \rho) = \tau_n(2\tau_n - \rho_n) + (\tau - \rho)^2(\lambda^2 + \tau\rho).$$

Let \mathcal{D} be the distribution on U defined by

$$\mathcal{D}(q) = \{X \in T_q(M) : \omega^n(X) = 0\}$$

for each point $q \in U$. Since

$$d\omega^n = \omega^a \wedge \omega_a^n = 0,$$

\mathcal{D} is completely integrable. Therefore at each point q of M there exists a local coordinate (y^1, \dots, y^n) such that each slice $y^n = \text{constant}$, say t , is an integral manifold of \mathcal{D} . We denote by $N(t)$ this slice. Restricting the orthonormal frame field $e_1, \dots, e_n, e_{n+1}, e_{n+2}$ to the slice $N(t)$, $\{e_a\}_{a=1, \dots, n-1}$ can be regarded as C^∞ orthonormal vector fields tangent to $N(t)$ and e_n, e_{n+1}, e_{n+2} can be also regarded as C^∞ orthonormal normal vector fields on $N(t)$ in E^{n+1} . Furthermore, restricting the dual frame ω^A and the connection form ω_B^A to the slice $N(t)$, we denote them by the same notation $\omega^i, \omega^a, \omega_j^i, \dots$ as those of the submanifold M in E^{n+2} . Then by (3.1) and Lemma 3.2, we see

$$\omega_a^n = \frac{\tau_n}{\tau - \rho} \omega^a,$$

$$\omega_a^{n+1} = \lambda \omega^a,$$

$$\omega_a^{n+2} = \tau \omega^a,$$

which mean that $N(t)$ is totally umbilic in E^{n+2} . Hence there exists an n -dimensional plane $L(t)$ in E^{n+2} such that $N(t)$ is a part of an $(n-1)$ -dimensional sphere in $L(t)$.

We denote by X the position vector of E^{n+2} and by H the mean curvature vector of $N(t)$ in E^{n+2} . It is then seen that the vector field

$$X + \frac{H}{\|H\|^2},$$

where we take X an arbitrary point of $N(t)$, is a constant vector in E^{n+2} and lies in $L(t)$, so it is called the *center* of $N(t)$. Moreover the tangent

space of $L(t)$ is spanned by e_1, \dots, e_{n-1} and $\frac{H}{\|H\|^2}$. And the tangent space of $N(t)$ is spanned by e_1, \dots, e_{n-1} . Precisely describing, the mean curvature vector H is expressed as

$$H = \frac{\tau_n}{\tau - \rho} e_n + \lambda e_{n+1} + \tau e_{n+2}.$$

Hence

$$\|H\|^2 = \frac{\tau_n^2 + (\lambda^2 + \tau^2)(\tau - \rho)^2}{(\tau - \rho)^2}.$$

Since $\tau \neq 0$ and $\tau - \rho \neq 0$, we have

$$\|H\| > 0.$$

We set

$$F = \tau_n^2 + (\lambda^2 + \tau^2)(\tau - \rho)^2.$$

Then

$$\frac{H}{\|H\|^2} = \frac{1}{F} \{ (\tau - \rho)\tau_n e_n + \lambda(\tau - \rho)^2 e_{n+1} + \tau(\tau - \rho)^2 e_{n+2} \}$$

Next we study the curve consisting of the centers of the slices. This curve is parametrized by t . We denote by $\bar{c}(t)$ this curve.

LEMMA 3.3. $\bar{c}(t)$ is orthogonal to e_1, \dots, e_{n-1} , $\frac{H}{\|H\|^2}$.

PROOF. Fix an arbitrary point q of M . Let $\tilde{c}(s)$ be the integral curve of e_n with the initial point q . We have only to show that

$$X + \frac{H}{\|H\|^2} \quad \left(\text{i. e. } \tilde{c}(s) + \frac{H}{\|H\|^2} \Big|_{\tilde{c}(s)} \right)$$

is orthogonal to e_1, \dots, e_{n-1} and $\frac{H}{\|H\|^2}$ when X moves on $\tilde{c}(s)$, because $\bar{c}(t)$

is the same curve as $\tilde{c}(s) + \frac{H}{\|H\|^2} \Big|_{\tilde{c}(s)}$ if we do not care the difference of the parameter.

Let (x^1, \dots, x^{n+2}) denote the natural coordinate of E^{n+2} . We set

$$\tilde{c}(s) = (\tilde{c}^1(s), \dots, \tilde{c}^{n+2}(s)),$$

$$\frac{H}{\|H\|^2} \Big|_{\tilde{c}(s)} = (h^1(s), \dots, h^{n+2}(s)),$$

in E^{n+2} . Then,

$$\frac{d}{ds} \left(\tilde{c}(s) + \frac{H}{\|H\|^2} \Big|_{\tilde{c}(s)} \right) = \left(\frac{d(\tilde{c}^1(s) + h^1(s))}{ds}, \dots, \frac{d(\tilde{c}^{n+2}(s) + h^{n+2}(s))}{ds} \right)$$

$$= D_{e_n} \left(X + \frac{H}{\|H\|^2} \right),$$

where D denotes the connection of E^{n+2} . It is sufficient to show that $D_{e_n} \left(X + \frac{H}{\|H\|^2} \right)$ is orthogonal to e_1, \dots, e_{n-1} and $\frac{H}{\|H\|^2}$.

Calculating straightforwardly, we get

$$D_{e_n} F = 2\tau_{nn}\tau_n + 2\tau\tau_n(\tau - \rho)^2 + 2(\lambda^2 + \tau^2)(\tau - \rho)(\tau_n - \rho_n),$$

where

$$F = \tau_n^2 + (\lambda^2 + \tau^2)(\tau - \rho)^2.$$

By (3.8) we obtain

$$(3.9) \quad (\tau - \rho)D_{e_n}(F) = 2F(2\tau_n - \rho_n).$$

Using (3.8) and (3.9), we get

$$(3.10) \quad D_{e_n} \frac{H}{\|H\|^2} = -\frac{1}{F} \{ \tau_n^2 e_n + \lambda(\tau - \rho)\tau_n e_{n+1} + (\tau - \rho)\tau\tau_n e_{n+2} \}.$$

Therefore

$$\begin{aligned} D_{e_n} \left(X + \frac{H}{\|H\|^2} \right) &= e_n + D_{e_n} \left(\frac{H}{\|H\|^2} \right) \\ &= \frac{1}{F} \{ (\lambda^2 + \tau^2)(\tau - \rho)^2 e_n \\ &\quad - \lambda(\tau - \rho)\tau_n e_{n+1} \\ &\quad - (\tau - \rho)\tau\tau_n e_{n+2} \}. \end{aligned}$$

This implies

$$\begin{aligned} \left\langle D_{e_n} \left(X + \frac{H}{\|H\|^2} \right), \frac{H}{\|H\|^2} \right\rangle &= 0, \\ \left\langle D_{e_n} \left(X + \frac{H}{\|H\|^2} \right), e_a \right\rangle &= 0. \end{aligned}$$

for $a=1, \dots, n-1$.

Q. E. D.

LEMMA 3.4. Let $\tilde{\mathcal{D}}$ be the distribution spanned by e_1, \dots, e_{n-1} and $\frac{H}{\|H\|^2}$. Then $\tilde{\mathcal{D}}$ is parallel on M .

PROOF. By the straightforward computation, we obtain

$$\begin{aligned} D_{e_a}(e_b) &= \omega_b^i(e_a)e_c + \frac{F}{(\tau - \rho)^2} \delta_{ab} \frac{H}{\|H\|^2} \in \tilde{\mathcal{D}}, \\ D_{e_n}(e_b) &= \omega_b^i(e_n)e_c \in \tilde{\mathcal{D}}, \\ D_{e_a} \left(\frac{H}{\|H\|^2} \right) &= D_{e_a} \left(X + \frac{H}{\|H\|^2} \right) - D_{e_a} X = -e_a \in \tilde{\mathcal{D}}, \end{aligned}$$

$$D_{e_n}\left(\frac{H}{\|H\|^2}\right) = -\frac{\tau_n}{\tau - \rho} \frac{H}{\|H\|^2} \in \tilde{\mathcal{D}}. \quad \text{Q. E. D.}$$

PROPOSITION 3.5. *Let the same assumption as Theorem 2.1 be satisfied. Suppose the condition II in Theorem 2.1 is satisfied on U and U has no point of constant sectional curvature. Then the shape of U becomes locally as follows:*

Let (x^1, \dots, x^{n+2}) denote the natural coordinate of E^{n+2} . There exist a curve $(x^1(t), x^2(t), 0, \dots, 0)$ in E^{n+2} and a C^∞ -function $\gamma(t) > 0$ such that U is represented locally by

$$(x^1(t), x^2(t), x^3, \dots, x^{n+2}),$$

where

$$(x^3)^2 + \dots + (x^{n+2})^2 = \gamma(t)^2.$$

PROOF. Lemma 3.3 implies that the curve $\bar{c}(t)$, which consists of the centers of the slices, is orthogonal to e_1, \dots, e_{n-1} and $\frac{H}{\|H\|^2}$. Furthermore, Lemma 3.4 implies that the distribution $\tilde{\mathcal{D}}$, which spanned by e_1, \dots, e_{n-1} and $\frac{H}{\|H\|^2}$, is parallel on M . From these two facts, it follows that $\bar{c}(t)$ must be a plane curve. Let this plane be x^1x^2 -plane, and let $(x^1(t), x^2(t), 0, \dots, 0)$ denote the curve $\bar{c}(t)$. On the other hand, each slice $N(t)$ is a part of $(n-1)$ -dimensional sphere of an n -dimensional plane $L(t)$ in E^{n+2} . Moreover the tangent space of $L(t)$ is spanned by e_1, \dots, e_{n-1} and $\frac{H}{\|H\|^2}$. Therefore each $L(t)$ is parallel mutually since $\tilde{\mathcal{D}}$ is parallel on U . Clearly the x^1x^2 -plane and $L(t)$ are orthogonal. Therefore we may take $x^3 \dots x^{n+2}$ -space to be parallel with each $L(t)$. Then U is represented as the situation mentioned above. Q. E. D.

4. The case of I in Theorem 2.1. In this section we shall study the shape of U on which the condition I in Theorem 2.1 is satisfied.

In this case there exist C^∞ orthonormal normal vector fields e_{n+1}, e_{n+2} and a C^∞ orthonormal frame field e_1, \dots, e_n on U such that

$$(4.1) \quad \begin{aligned} \omega_s^{n+1} &= \lambda \omega^s, & \omega_s^{n+2} &= \tau \omega^s, \\ \omega_{n-1}^{n+1} &= \mu \omega^{n-1}, & \omega_{n-1}^{n+2} &= \tau \omega^{n-1}, \\ \omega_n^{n+1} &= \lambda \omega^n, & \omega_n^{n+2} &= \rho \omega^n, \end{aligned}$$

where λ, τ, μ and ρ are C^∞ -functions on U and $\lambda \neq \mu, \tau \neq \rho$ on U .

In this section indices A, B, \dots run over the range $\{1, \dots, n+2\}$ and

s, t, \dots move from 1 to $n-2$. For indices s, t, \dots , the Einstein's convention is used.

LEMMA 4.1. *A point q in M is of constant sectional curvature if and only if $\lambda=0$ and $\tau=0$ at q .*

Using the same method as Lemma 3.1, we will soon prove Lemma 4.1.

Since we have assumed that U has no point of constant sectional curvature, we may set $\lambda \neq 0$ or $\tau \neq 0$ on U . Now we assume $\tau \neq 0$.

We set

$$\begin{aligned}d\lambda &= \lambda_s \omega^s + \lambda_{n-1} \omega^{n-1} + \lambda_n \omega^n, \\d\tau &= \tau_s \omega^s + \tau_{n-1} \omega^{n-1} + \tau_n \omega^n, \\d\mu &= \mu_s \omega^s + \mu_{n-1} \omega^{n-1} + \mu_n \omega^n, \\d\rho &= \rho_s \omega^s + \rho_{n-1} \omega^{n-1} + \rho_n \omega^n.\end{aligned}$$

Using the same method getting Lemma 3.2, we obtain Lemma 4.2.

LEMMA 4.2.

$$\begin{aligned}(4.3) \quad \omega_{n-1}^s &= -\frac{\lambda_{n-1} + \tau C_{n-1}}{\lambda - \mu} \omega^s, \\ \omega_n^s &= -\frac{\lambda \lambda_n + \tau \tau_n}{\tau(\tau - \rho)} \omega^s, \\ \omega_{n+2}^{n+1} &= C_{n-1} \omega^{n-1} - \frac{\lambda_n}{\tau} \omega^n, \\ (4.4) \quad \omega_n^{n-1} &= \frac{-\lambda_n + \mu_n}{\lambda - \mu} \omega^{n-1} + \frac{\tau_{n-1} - \rho_{n-1}}{\tau - \rho} \omega^n, \\ (4.5) \quad \lambda_s &= 0, \quad \tau_s = 0, \quad \mu_s = 0, \quad \rho_s = 0, \\ &\lambda C_{n-1} - \tau_{n-1} = 0,\end{aligned}$$

where we set

$$\omega_{n+2}^{n+1}(e_{n-1}) = C_{n-1}.$$

Now, we set

$$\begin{aligned}d\lambda_{n-1} &= \lambda_{n-1,t} \omega^t + \lambda_{n-1,n-1} \omega^{n-1} + \lambda_{n-1,n} \omega^n, \\d\lambda_n &= \lambda_{n,t} \omega^t + \lambda_{n,n-1} \omega^{n-1} + \lambda_{n,n} \omega^n, \\d\tau_{n-1} &= \tau_{n-1,t} \omega^t + \tau_{n-1,n-1} \omega^{n-1} + \tau_{n-1,n} \omega^n, \\d\tau_n &= \tau_{n,t} \omega^t + \tau_{n,n-1} \omega^{n-1} + \tau_{n,n} \omega^n, \\d\mu_{n-1} &= \mu_{n-1,t} \omega^t + \mu_{n-1,n-1} \omega^{n-1} + \mu_{n-1,n} \omega^n,\end{aligned}$$

$$\begin{aligned}
d\mu_n &= \mu_{n,t}\omega^t + \mu_{n,n-1}\omega^{n-1} + \mu_{n,n}\omega^n, \\
d\rho_{n-1} &= \rho_{n-1,t}\omega^t + \rho_{n-1,n-1}\omega^{n-1} + \rho_{n-1,n}\omega^n, \\
d\rho_n &= \rho_{n,t}\omega^t + \rho_{n,n-1}\omega^{n-1} + \rho_{n,n}\omega^n.
\end{aligned}$$

Because of $d^2=0$, we obtain

$$\begin{aligned}
(4.7) \quad & \lambda_{n-1,t}=0, \quad \lambda_{n,t}=0, \quad \tau_{n-1,t}=0, \quad \tau_{n,t}=0, \\
& \mu_{n-1,t}=0, \quad \mu_{n,t}=0, \quad \rho_{n-1,t}=0, \quad \rho_{n,t}=0.
\end{aligned}$$

We set

$$dC_{n-1}(e_s) = C_{n-1,s}.$$

By taking the exterior derivative of (4.3), we get

$$(4.8) \quad C_{n-1,s} = 0,$$

$$\begin{aligned}
(4.9) \quad & (\lambda - \mu)\tau(\tau - \rho)\{e_{n-1}(\lambda_{n-1} + \tau C_{n-1})\} \\
& = \tau(\tau - \rho)(\lambda_{n-1} - \mu_{n-1})(\lambda_{n-1} + \tau C_{n-1}) \\
& \quad - (\lambda - \mu)(-\lambda_n + \mu_n)(\lambda\lambda_n + \tau\tau_n) \\
& \quad + (\lambda - \mu)^2\tau(\tau - \rho)(\lambda\mu + \tau^2) \\
& \quad + \tau(\tau - \rho)(\lambda_{n-1} + \tau C_{n-1})^2, \\
(4.10) \quad & (\lambda - \mu)\tau(\tau - \rho)^2\{e_n(\lambda_{n-1} + \tau C_{n-1})\} \\
& = \tau(\tau - \rho)^2(\lambda_n - \mu_n)(\lambda_{n-1} + \tau C_{n-1}) \\
& \quad - (\lambda - \mu)^2(\lambda\lambda_n + \tau\tau_n)(\tau_{n-1} - \rho_{n-1}) \\
& \quad + (\lambda - \mu)(\tau - \rho)(\lambda_{n-1} + \tau C_{n-1})(\lambda\lambda_n + \tau\tau_n),
\end{aligned}$$

$$\begin{aligned}
(4.11) \quad & \tau(\tau - \rho)(\lambda - \mu)^2\{e_{n-1}(\lambda\lambda_n + \tau\tau_n)\} \\
& = (\lambda - \mu)^2(\lambda\lambda_n + \tau\tau_n)\{\tau_{n-1}(\tau - \rho) + \tau(\tau_{n-1} - \rho_{n-1})\} \\
& \quad + \tau^2(\tau - \rho)^2(-\lambda_n + \mu_n)(\lambda_{n-1} + \tau C_{n-1}) \\
& \quad + \tau(\tau - \rho)(\lambda - \mu)(\lambda\lambda_n + \tau\tau_n)(\lambda_{n-1} + \tau C_{n-1}),
\end{aligned}$$

$$\begin{aligned}
(4.12) \quad & \tau(\tau - \rho)(\lambda - \mu)\{e_n(\lambda\lambda_n + \tau\tau_n)\} \\
& = (\lambda - \mu)(\lambda\lambda_n + \tau\tau_n)\{\tau_n(\tau - \rho) + \tau(\tau_n - \rho_n)\} \\
& \quad + \tau^2(\tau - \rho)(\tau_{n-1} - \rho_{n-1})(\lambda_{n-1} + \tau C_{n-1}) \\
& \quad + \tau^2(\tau - \rho)^2(\lambda - \mu)(\lambda^2 + \rho\tau) \\
& \quad + (\lambda\lambda_n + \tau\tau_n)^2(\lambda - \mu).
\end{aligned}$$

Let \mathcal{D}_1 be the distribution on U defined by

$$\mathcal{D}_1(q) = \{X \in T_q(M) : \omega^{n-1}(X) = 0, \omega^n(X) = 0\}.$$

Since $d\omega^{n-1} = \omega^n \wedge \omega_n^{n-1}$ and $d\omega^n = \omega^{n-1} \wedge \omega_n^n$, \mathcal{D}_1 is completely integrable.

On the other hand, let \mathcal{D}_2 be the distribution on U defined by

$$\mathcal{D}_2(q) = \{X \in T_q(M) : \omega^s(X) = 0 \ (s=1, \dots, n-2)\}.$$

(4.3) implies that \mathcal{D}_2 is completely integrable.

From above two facts it is obtained that at each point q on M there exists a local coordinate (y^1, \dots, y^n) such that the slice $y^{n-1} = \text{constant}$, $y^n = \text{constant}$ is an integral manifold of \mathcal{D}_1 and the slice $y^s = \text{constant}$ ($s=1, \dots, n-2$) is an integral manifold of \mathcal{D}_2 . Now, we denote by $N(u, v)$ the slice $y^{n-1} = \text{constant} = u$ and $y^n = \text{constant} = v$, and denote by S the slice $y^s = 0$ ($s=1, \dots, n-2$).

Restricting the orthogonal frame field e_1, \dots, e_n, e_{n+1} and e_{n+2} to the slice $N(u, v)$, $\{e_s\}_{s=1, \dots, n-2}$ can be regarded as C^∞ orthogonal vector fields tangent to $N(u, v)$ and $e_{n-1}, e_n, e_{n+1}, e_{n+2}$ can be also regarded as C^∞ -orthonormal normal vector fields on $N(u, v)$ in E^{n+2} . Furthermore, restricting the dual frame ω^A and the connection form ω_A^B to the slice $N(u, v)$, we denote them by the same notation $\omega^i, \omega^\alpha, \omega_j^i, \dots$ as those of the submanifold M in E^{n+2} . By (4.1) and Lemma 4.2, we see

$$\omega_s^{n-1} = \frac{\lambda_{n-1} + \tau C_{n-1}}{\lambda - \mu} \omega^s,$$

$$\omega_s^n = \frac{\lambda \lambda_n + \tau \tau_n}{\tau(\tau - \rho)} \omega^s,$$

$$\omega_s^{n+2} = \lambda \omega^s,$$

$$\omega_s^{n+2} = \tau \omega^s,$$

which mean that $N(u, v)$ is totally umbilic in E^{n+2} . Hence there exists $(n-1)$ -dimensional plane $L(u, v)$ in E^{n+2} such that $N(u, v)$ is a part of an $(n-2)$ -dimensional sphere in $L(u, v)$. Let X denote the position vector of E^{n+2} , and H denote the mean curvature vector of $N(u, v)$ in E^{n+2} . It is seen that the vector field

$$X + \frac{H}{\|H\|^2},$$

where we take X an arbitrary point of $N(u, v)$, is a constant vector in E^{n+2} and lies in $L(u, v)$, so it is called the *center* of $N(u, v)$. Moreover the tangent space of $L(u, v)$ is spanned by e_1, \dots, e_{n-2} and $\frac{H}{\|H\|^2}$.

LEMMA 4.3. Let \mathcal{D}_3 be the distribution spanned by e_1, \dots, e_{n-2} and

$\frac{H}{\|H\|^2}$. Then $D_{e_{n-1}}\left(X + \frac{H}{\|H\|^2}\right)$ and $D_{e_n}\left(X + \frac{H}{\|H\|^2}\right)$ are orthogonal to \mathcal{D}_3 .

PROOF. Precisely describing,

$$H = \frac{\lambda_{n-1} + \tau C_{n-1}}{\lambda - \mu} e_{n-1} + \frac{\lambda \lambda_n + \tau \tau_n}{\tau(\tau - \rho)} e_n + \lambda e_{n+1} + \tau e_{n+2}.$$

Since $\tau \neq 0$, $\tau \neq \rho$ and $\lambda \neq \mu$, we get

$$\|H\|^2 > 0.$$

We set

$$F = (\lambda_{n-1} + \tau C_{n-1})^2 \tau^2 (\tau - \rho)^2 + (\lambda \lambda_n + \tau \tau_n)^2 (\lambda - \mu)^2 \\ + (\lambda^2 + \tau^2) (\lambda - \mu)^2 \tau^2 (\tau - \rho)^2.$$

Then

$$\frac{H}{\|H\|^2} = \frac{1}{F} \{ (\lambda - \mu) \tau^2 (\tau - \rho)^2 (\lambda_{n-1} + \tau C_{n-1}) e_{n-1} \\ + (\lambda - \mu)^2 \tau (\tau - \rho) (\lambda \lambda_n + \tau \tau_n) e_n \\ + \lambda (\lambda - \mu)^2 \tau^2 (\tau - \rho)^2 e_{n+1} \\ + (\lambda - \mu)^2 \tau^3 (\tau - \rho)^2 e_{n+2} \}.$$

From (4.5), (4.9) and (4.11), we get

$$(4.13) \quad \tau(\tau - \rho)(\lambda - \mu) D_{e_{n-1}} F = 2F \{ \tau(\tau - \rho)(\lambda_{n-1} + \tau C_{n-1}) \\ + \tau_{n-1}(\tau - \rho)(\lambda - \mu) \\ + \tau(\tau_{n-1} - \rho_{n-1})(\lambda - \mu) \\ + \tau(\tau - \rho)(\lambda_{n-1} - \mu_{n-1}) \}.$$

By (4.13), (4.9) and (4.13), we get

$$(4.14) \quad D_{e_{n-1}} \frac{H}{\|H\|^2} \\ = -\frac{1}{F} \{ \tau^2 (\tau - \rho)^2 (\lambda_{n-1} + \tau C_{n-1})^2 e_{n-1} \\ + (\lambda - \mu) \tau (\tau - \rho) (\lambda \lambda_n + \tau \tau_n) (\lambda \lambda_{n-1} + \tau \tau_{n-1}) e_n \\ + \lambda (\lambda - \mu) \tau^2 (\tau - \rho)^2 (\lambda_{n-1} + \tau C_{n-1}) e_{n+1} \\ + (\lambda - \mu) \tau^3 (\tau - \rho)^2 (\lambda_{n-1} + \tau C_{n-1}) e_{n+2} \}.$$

Hence

$$\begin{aligned}
 (4.15) \quad & D_{e_{n-1}}\left(X + \frac{H}{\|H\|^2}\right) \\
 &= e_{n-1} + D_{e_{n-1}}\frac{H}{\|H\|^2} \\
 &= \frac{1}{F}\{(\lambda - \mu)^2(\lambda\lambda_n + \tau\tau_n)^3 \\
 &\quad + (\lambda^2 + \tau^2)(\lambda - \mu)^2\tau^2(\tau - \rho)^2\}e_{n-1} \\
 &\quad - (\lambda - \mu)\tau(\tau - \rho)(\lambda\lambda_n + \tau\tau_n)(\lambda_{n-1} + \tau C_{n-1})e_n \\
 &\quad - \lambda(\lambda - \mu)\tau^2(\tau - \rho)^2(\lambda_{n-1} + \tau C_{n-1})e_{n+1} \\
 &\quad - (\lambda - \mu)\tau^3(\tau - \rho)^2(\lambda_{n-1} + \tau C_{n-1})e_{n+2}.
 \end{aligned}$$

This implies

$$\begin{aligned}
 \left\langle D_{e_{n-1}}\left(X + \frac{H}{\|H\|^2}\right), \frac{H}{\|H\|^2} \right\rangle &= 0, \\
 \left\langle D_{e_{n-1}}\left(X + \frac{H}{\|H\|^2}\right), e_s \right\rangle &= 0,
 \end{aligned}$$

where $s=1, \dots, n-2$. Therefore $D_{e_{n-1}}\left(X + \frac{H}{\|H\|^2}\right)$ is orthogonal to \mathcal{D}_3 .

By (4.10) and (4.12), we get

$$\begin{aligned}
 \tau(\tau - \rho)(\lambda - \mu)D_{e_n}F &= 2F\{\tau_n(\tau - \rho)(\lambda - \mu) + (\tau_n - \rho_n)\tau(\lambda - \mu) \\
 &\quad + (\lambda_n - \mu_n)\tau(\tau - \rho) + (\lambda\lambda_n + \tau\tau_n)(\lambda - \mu)\}.
 \end{aligned}$$

Therefore we get

$$\begin{aligned}
 & D_{e_n}\left(X + \frac{H}{\|H\|^2}\right) \\
 &= \frac{1}{F}\{-(\lambda - \mu)\tau(\tau - \rho)(\lambda_{n-1} + \tau C_{n-1})(\lambda\lambda_n + \tau\tau_n)e_{n-1} \\
 &\quad + \{\tau^2(\tau - \rho)^2(\lambda_{n-1} + \tau C_{n-1})^2 \\
 &\quad + (\lambda^2 + \tau^2)(\lambda - \mu)^2\tau^2(\tau - \rho)^2\}e_n \\
 &\quad - \lambda(\lambda - \mu)^2\tau(\tau - \rho)(\lambda\lambda_n + \tau\tau_n)e_{n+1} \\
 &\quad - (\lambda - \mu)^3\tau^2(\tau - \rho)(\lambda\lambda_n + \tau\tau_n)e_{n+2}.
 \end{aligned}$$

This implies

$$\begin{aligned}
 \left\langle D_{e_n}\left(X + \frac{H}{\|H\|^2}\right), \frac{H}{\|H\|^2} \right\rangle &= 0, \\
 \left\langle D_{e_n}\left(X + \frac{H}{\|H\|^2}\right), e_s \right\rangle &= 0,
 \end{aligned}$$

where $s=1, \dots, n-2$. Therefore $D_{e_n}\left(X+\frac{H}{\|H\|^2}\right)$ is orthogonal to \mathcal{D}_3 .

Q. E. D.

LEMMA 4.4. \mathcal{D}_3 , which is spanned by e_1, \dots, e_{n-2} and $\frac{H}{\|H\|^2}$ is parallel on M .

PROOF. Straightforwardly computing, we get

$$D_{e_s}e_t = \omega_t^r(e_s)e_r + \frac{\delta_{st}F}{(\lambda-\mu)^2\tau^2(\tau-\rho)^2} \frac{H}{\|H\|^2} \in \mathcal{D}_3,$$

$$D_{e_{n-1}}e_t = \omega_t^s(e_{n-1})e_s \in \mathcal{D}_3,$$

$$D_{e_n}e_t = \omega_t^s(e_n)e_s \in \mathcal{D}_3,$$

$$D_{e_s}\frac{H}{\|H\|^2} = D_{e_s}\left(X+\frac{H}{\|H\|^2}\right) - e_s = -e_s \in \mathcal{D}_3.$$

$$D_{e_{n-1}}\frac{H}{\|H\|^2} = -\frac{\lambda_{n-1}+\tau C_{n-1}}{\lambda-\mu} \frac{H}{\|H\|^2} \in \mathcal{D}_3.$$

$$D_{e_n}\frac{H}{\|H\|^2} = -\frac{\lambda\lambda_n+\tau\tau_n}{\tau(\tau-\rho)} \frac{H}{\|H\|^2} \in \mathcal{D}_3.$$

Q. E. D.

LEMMA 4.5. Let Ψ be a map from S to E^{n+2} defined by $\Psi(q) = X_q + \frac{H}{\|H\|^2}\Big|_q$ for $q \in S$, where S is the slice defined by $y^s=0$ ($s=1, \dots, n-2$). Then Ψ is an immersion.

PROOF. We have only to show that $d\Psi(e_{n-1})$ and $d\Psi(e_n)$ are linearly independent. Now,

$$\begin{aligned} d\Psi(e_{n-1}) &= D_{e_{n-1}}\left(X+\frac{H}{\|H\|^2}\right) \\ &= \frac{1}{F}(\{\lambda\lambda_n+\tau\tau_n\}(\lambda-\mu)^2 \\ &\quad + (\lambda^2+\tau^2)(\lambda-\mu)^2\tau^2(\tau-\rho)^2)e_{n-1} \\ &\quad - (\lambda-\mu)\tau(\tau-\rho)(\lambda\lambda_n+\tau\tau_n)(\lambda_{n-1}+\tau C_{n-1})e_n \\ &\quad - \lambda(\lambda-\mu)\tau^2(\tau-\rho)^2(\lambda_{n-1}+\tau C_{n-1})e_{n+1} \\ &\quad - (\lambda-\mu)\tau^3(\tau-\rho)^2(\lambda_{n-1}+\tau C_{n-1})e_{n+2}. \end{aligned}$$

$$\begin{aligned} d\Psi(e_n) &= \frac{1}{F}(-(\lambda-\mu)\tau(\tau-\rho)(\lambda_{n-1}+\tau C_{n-1})(\lambda\lambda_n+\tau\tau_n)e_{n-1} \\ &\quad + \{\tau^2(\tau-\rho)^2(\lambda_{n-1}+\tau C_{n-1})\}^2) \end{aligned}$$

$$\begin{aligned}
 &+(\lambda^2 + \tau^2)(\lambda - \mu)^2 \tau^2 (\tau - \rho)^2 e_n \\
 &-\lambda(\lambda - \mu)^2 \tau (\tau - \rho)(\lambda \lambda_n + \tau \tau_n) e_{n+1} \\
 &-(\lambda - \mu)^2 \tau^2 (\tau - \rho)(\lambda \lambda_n + \tau \tau_n) e_{n+2}.
 \end{aligned}$$

Since $\lambda \neq \mu$, $\tau \neq \rho$ and $\tau \neq 0$, it is clear that $d\Psi(e_{n-1})$ and $d\Psi(e_n)$ are linearly independent. Q. E. D.

Lemma 4.5 implies that $\Psi(S) = \left\{ X + \frac{H}{\|H\|^2} : X \in S \right\}$, which consists of the centers, becomes a 2-dimensional submanifold of E^{n+2} .

PROPOSITION 4.6. *Let the same assumption as Theorem 2.1 be satisfied. Suppose the condition I in Theorem 2.1 is satisfied on U and U has no point of constant sectional curvature. Then the shape of U becomes as follows:*

Let (x^1, \dots, x^{n+2}) denote the natural coordinate of E^{n+2} . There exist a surface $(x^1(u, v), x^2(u, v), x^3(u, v), 0, \dots, 0)$ in $x^1 x^2 x^3$ -space in E^{n+2} and a C^∞ -function $\gamma(u, v) (> 0)$ such that U is locally represented by

$$(x^1(u, v), x^2(u, v), x^3(u, v), x^4, \dots, x^{n+2}),$$

where

$$(x^4)^2 + \dots + (x^{n+2})^2 = \gamma(u, v)^2.$$

PROOF. The tangent space of $\Psi(S)$ is spanned by $D_{e_{n-1}}\left(X + \frac{H}{\|H\|^2}\right)$ and $D_{e_n}\left(X + \frac{H}{\|H\|^2}\right)$. On the other hand, Lemma 4.3 shows that both $D_{e_{n-1}}\left(X + \frac{H}{\|H\|^2}\right)$ and $D_{e_n}\left(X + \frac{H}{\|H\|^2}\right)$ are orthogonal to e_1, \dots, e_{n-2} and $\frac{H}{\|H\|^2}$. Moreover Lemma 4.4 implies that the distribution spanned by e_1, \dots, e_{n-2} and $\frac{H}{\|H\|^2}$ is parallel on M. Therefore $D_{e_{n-1}}\left(X + \frac{H}{\|H\|^2}\right)$ and $D_{e_n}\left(X + \frac{H}{\|H\|^2}\right)$ are contained in a fixed 3-dimensional space. Let this space be $x^1 x^2 x^3$ -space and let $(x^1(u, v), x^2(u, v), x^3(u, v), 0, \dots, 0)$ represent the surface $\Psi(S)$. It has been shown that each slice $N(u, v)$ is a part of an $(n-2)$ -dimensional sphere in $L(u, v)$, which is an $(n-1)$ -dimensional plane in E^{n+2} . Since the tangent space of $L(u, v)$ is spanned by e_1, \dots, e_{n-2} , $\frac{H}{\|H\|^2}$ and the distribution which is spanned by e_1, \dots, e_{n-2} , $\frac{H}{\|H\|^2}$ is parallel on M, each $L(u, v)$ is parallel. Clearly the $x^1 x^2 x^3$ -space is orthogonal to $L(u, v)$. Therefore we may choose $x^4 \dots x^{n+2}$ -space to be parallel

with $L(u, v)$. Then U is represented as the situation mentioned above.

Q. E. D.

It is easy to show that the Theorem 2.1, Proposition 3.5 and Proposition 4.6 lead the Theorem in the introduction.

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