# $L_{k}$-2-TYPE HYPERSURFACES IN HYPERBOLIC SPACES 

Pascual Lucas* and Héctor-Fabián Ramírez-Ospina


#### Abstract

In this article, we study $L_{k}$-finite-type hypersurfaces $M^{n}$ of a hyperbolic space $\mathbb{H}^{n+1} \subset \mathbb{R}_{1}^{n+2}$, for $k \geq 1$. In the 3-dimensional case, we obtain the following classification result. Let $\psi: M^{3} \rightarrow \mathbb{H}^{4} \subset \mathbb{R}_{1}^{5}$ be an orientable hypersurface with constant $k$-th mean curvature $H_{k}$, which is not totally umbilical. Then $M^{3}$ is of $L_{k}$-2-type if and only if $M^{3}$ is an open portion of a standard Riemannian product $\mathbb{H}^{1}\left(r_{1}\right) \times \mathbb{S}^{2}\left(r_{2}\right)$ or $\mathbb{H}^{2}\left(r_{1}\right) \times \mathbb{S}^{1}\left(r_{2}\right)$, with $-r_{1}^{2}+r_{2}^{2}=-1$. In the $n$-dimensional case, we show that a hypersurface $M^{n} \subset \mathbb{H}^{n+1}$, with constant $k$-th mean curvature $H_{k}$ and having at most two distinct principal curvatures, is of $L_{k}$-2-type if and only if $M^{n}$ is an open portion of a Riemannian product $\mathbb{H}^{m}\left(r_{1}\right) \times \mathbb{S}^{n-m}\left(r_{2}\right)$, with $-r_{1}^{2}+r_{2}^{2}=-1$, for some integer $m \in\{1, \ldots, n-1\}$. In the case $k=n-1$ we drop the condition on the principal curvatures of the hypersurface $M^{n}$, and prove that if $M^{n} \subset \mathbb{H}^{n+1}$ is an orientable $H_{n-1}$-hypersurface of $L_{n-1}-2$-type then its Gauss-Kronecker curvature $H_{n}$ is a nonzero constant.


## 1. Introduction

Submanifolds of finite type were introduced by B.Y. Chen, whose first results were gathered in his book [7] (see also [8]). Although the first definition was given for a compact submanifold in the Euclidean space, Chen extended the concept to noncompact submanifolds in Euclidean or pseudo-Euclidean spaces, [9, 10]. A detailed survey of the results on this subject, up to 1996, was given by Chen in [14], and in a recent article [15], the author provides a detailed account of recent development on the problems and conjectures about finite type submanifolds.

The Laplacian operator $\Delta$ can be seen as the first one of a sequence of $n$ operators $L_{0}=\Delta, L_{1}, \ldots, L_{n-1}$, where $L_{k}$ stands for the linearized operator of the first variation

[^0]of the $(k+1)$-th mean curvature arising from normal variations of the hypersurface (see, for instance, [24]). These operators $L_{k}$ are given by $L_{k}(f)=\operatorname{tr}\left(P_{k} \circ \nabla^{2} f\right)$ for a smooth function $f$ on $M$, where $P_{k}$ denotes the $k$-th Newton transformation associated to the second fundamental form of the hypersurface and $\nabla^{2} f$ denotes the self-adjoint linear operator metrically equivalent to the Hessian of $f$.

As an extension of finite type theory, S.M.B. Kashani [17] introduced the notion of $L_{k}$-finite-type hypersurface in the Euclidean space. In general, a submanifold $M^{n}$ in $\mathbb{R}^{m}$ is said to be of $L_{k}$-finite-type if the position vector $\psi: M^{n} \rightarrow \mathbb{R}^{m}$ of $M^{n}$ into $\mathbb{R}^{m}$ admits the following finite spectral decomposition

$$
\psi=a+\psi_{1}+\cdots+\psi_{q}, \quad L_{k} \psi_{t}=\lambda_{t} \psi_{t},
$$

where $a$ is a constant vector, $\lambda_{t}$ are constants and $\psi_{t}$ are non-constant $\mathbb{R}^{m}$-valued maps on $M^{n}$. If all $\lambda_{t}$ 's are mutually different, $M^{n}$ is said to be of $L_{k}-q$-type, and if one of $\lambda_{t}$ is zero $M^{n}$ is said to be of $L_{k}$-null- $q$-type. Naturally, that definition is also valid for a pseudo-Riemannian submanifold $M_{t}^{n}$ into the pseudo-Euclidean space $\mathbb{R}_{s}^{m}$.

In [21], the authors, by using results from [1], show that $k$-minimal Euclidean hypersurfaces and open portions of hyperspheres are the only $L_{k}$-1-type hypersurfaces in $\mathbb{R}^{n+1}$. As for hypersurfaces of $L_{k}$-2-type in $\mathbb{R}^{n+1}$, the authors show that if $M^{n}$ is a hypersurface with at most two distinct principal curvatures, then (i) $M^{n}$ is not of $L_{n-1}$-null-2-type (Theorem 3.5); and (ii) $M^{n}$ is of $L_{k}$-null-2-type $(k \neq n-1)$ if and only if $M$ is locally isometric to a generalized cylinder (Theorems 3.11 and 3.12).

In [20], the authors study $L_{k}$-2-type hypersurfaces in a hypersphere $\mathbb{S}^{4} \subset \mathbb{R}^{5}$. Since the case $k=0$ corresponds to the classical one, which has been well studied (see, e.g., [11], [12] and [16], among others), the authors concentrate in cases $k=1$ and $k=2$, and show the following result:

Theorem A. Let $\psi: M^{3} \rightarrow \mathbb{S}^{4} \subset \mathbb{R}^{5}$ be an orientable $H_{k}$-hypersurface, which is not an open portion of a hypersphere. Then $M^{3}$ is of $L_{k}$-2-type if and only if $M^{3}$ is a Clifford tori $\mathbb{S}^{1}\left(r_{1}\right) \times \mathbb{S}^{2}\left(r_{2}\right)$, $r_{1}^{2}+r_{2}^{2}=1$, for appropriate radii, or a tube $T^{r}\left(V^{2}\right)$ of appropriate constant radius $r$ around the Veronese embedding $V^{2}$ of the real projective plane $\mathbb{R} P^{2}(\sqrt{3})$.

In this paper we extend this result to hypersurfaces in a hyperbolic space. The case $k=0$ was studied by Chen, [13], in the $n$-dimensional case. He proved (i) that every 2-type hypersurface of the hyperbolic space has nonzero constant mean curvature and constant scalar curvature, and (ii) that there exists no compact 2-type hypersurfaces in the hyperbolic space.

After a section devoted to preliminaries and basic results we proceed, in the third section, to compute some formulae required to present the examples. In section 4 we present the main results in dimension three, which we can gather in the following classification theorem:

Theorem B. Let $\psi: M^{3} \rightarrow \mathbb{H}^{4} \subset \mathbb{R}_{1}^{5}$ be an orientable $H_{k}$-hypersurface, which is not totally umbilical. Then $M^{3}$ is of $L_{k}$-2-type if and only if $M^{3}$ is a standard Riemannian product $\mathbb{H}^{1}\left(r_{1}\right) \times \mathbb{S}^{2}\left(r_{2}\right)$ or $\mathbb{H}^{2}\left(r_{1}\right) \times \mathbb{S}^{1}\left(r_{2}\right)$, with $-r_{1}^{2}+r_{2}^{2}=-1$.

In the final section, we extend the previous result to $n$-dimensional hypersurfaces in the hyperbolic space $\mathbb{H}^{n+1}$ as follows.

Theorem C. Let $\psi: M^{n} \rightarrow \mathbb{H}^{n+1} \subset \mathbb{R}_{1}^{n+2}$ be an orientable $H_{k}$-hypersurface and assume that $M^{n}$ has at most two distinct principal curvatures. Then $M^{n}$ is of $L_{k}$-2-type if and only if $M^{n}$ is an open portion of $\mathbb{H}^{m}\left(-\sqrt{1+r^{2}}\right) \times \mathbb{S}^{n-m}(r)$, for some positive integer $m, 1 \leq m \leq n-1$, and for some positive number $r$.

We wish to thank the referee for his/her comments and suggestions that have improved the original manuscript.

## 2. Preliminaries and Lemma

Let $\mathbb{R}_{1}^{5}$ be the 5-dimensional Lorentzian space with the standard flat metric $g$ given by

$$
g=-\mathrm{d} x_{1}^{2}+\sum_{i=2}^{5} \mathrm{~d} x_{i}^{2},
$$

where $\left(x_{1}, \ldots, x_{5}\right)$ is a rectangular coordinate system of $\mathbb{R}_{1}^{5}$. For a positive number $r$ and a point $c \in \mathbb{R}_{1}^{5}$ we denote by $\mathbb{H}^{4}(c,-r)$ the (connected) hyperbolic space centered at $c$ with radius $r$, which is embedded standardly in $\mathbb{R}_{1}^{5}$ by

$$
\mathbb{H}^{4}(c,-r)=\left\{x \in \mathbb{R}_{1}^{5} \mid\langle x-c, x-c\rangle=-r^{2}, \text { and } x_{1}>0\right\},
$$

where $\langle$,$\rangle denotes the Lorentzian inner product on \mathbb{R}_{1}^{5}$. To simplify the notation, we write $\mathbb{H}^{4}(-r) \equiv \mathbb{H}^{4}(0,-r)$ and $\mathbb{H}^{4} \equiv \mathbb{H}^{4}(0,-1)$. We will also use $\langle$,$\rangle to denote the$ flat metric $g$. Without loss of generality, we assume that $c=0$ and $r=1$.

Let $\psi: M^{3} \rightarrow \mathbb{H}^{4} \subset \mathbb{R}_{1}^{5}$ be an isometric immersion of a connected orientable hypersurface $M^{3}$ with Gauss map $N$. We denote by $\nabla^{0}$, $\bar{\nabla}$ and $\nabla$ the Levi-Civita connections on $\mathbb{R}_{1}^{5}, \mathbb{H}^{4}$ and $M^{3}$, respectively. Then the Gauss and Weingarten formulae are given by [22]

$$
\begin{aligned}
\nabla_{X}^{0} Y & =\nabla_{X} Y+\langle S X, Y\rangle N+\langle X, Y\rangle \psi \\
S X & =-\bar{\nabla}_{X} N=-\nabla_{X}^{0} N
\end{aligned}
$$

for all tangent vector fields $X, Y \in \mathfrak{X}\left(M^{3}\right)$, where $S: \mathfrak{X}\left(M^{3}\right) \longrightarrow \mathfrak{X}\left(M^{3}\right)$ stands for the shape operator (or Weingarten endomorphism) of $M^{3}$, with respect to the chosen orientation $N$.

As is well known, for every point $p \in M^{3}, S$ defines a linear self-adjoint endomorphism on the tangent space $T_{p} M^{3}$, and its eigenvalues $\kappa_{1}(p), \kappa_{2}(p)$ and $\kappa_{3}(p)$ are the principal curvatures of the hypersurface. The characteristic polynomial $Q_{S}(t)$ of $S$ is defined by

$$
Q_{S}(t)=\operatorname{det}(t I-S)=\left(t-\kappa_{1}\right)\left(t-\kappa_{2}\right)\left(t-\kappa_{3}\right)=t^{3}+a_{1} t^{2}+a_{2} t+a_{3},
$$

where the coefficients of $Q_{S}(t)$ are given by
$a_{1}=-\left(\kappa_{1}+\kappa_{2}+\kappa_{3}\right), \quad a_{2}=\kappa_{1} \kappa_{2}+\kappa_{1} \kappa_{3}+\kappa_{2} \kappa_{3}, \quad a_{3}=-\kappa_{1} \kappa_{2} \kappa_{3}$.
These coefficients can be expressed in terms of the traces of $S^{j}$ as follows:

$$
\begin{align*}
& a_{1}=-\operatorname{tr}(S), \\
& a_{2}=-\frac{1}{2} \operatorname{tr}\left(S^{2}\right)+\frac{1}{2} \operatorname{tr}(S)^{2},  \tag{1}\\
& a_{3}=-\frac{1}{3} \operatorname{tr}\left(S^{3}\right)+\frac{1}{2} \operatorname{tr}\left(S^{2}\right) \operatorname{tr}(S)-\frac{1}{6} \operatorname{tr}(S)^{3} .
\end{align*}
$$

The $k$-th mean curvature $H_{k}$ or mean curvature of order $k$ of $M^{3}$ in $\mathbb{H}^{4}$ is defined by

$$
\binom{3}{k} H_{k}=(-1)^{k} a_{k}, \quad \text { with } H_{0}=1
$$

We say that $M^{3}$ is an $H_{k}$-hypersurface if its $k$-th mean curvature $H_{k}$ is constant. If $H_{k+1}=0$, we then say that $M^{3}$ is a $k$-minimal hypersurface; a 0 -minimal hypersurface is nothing but a minimal hypersurface in $\mathbb{H}^{4}$.

The $k$-th Newton transformation of $M^{3}$ is the operator $P_{k}: \mathfrak{X}\left(M^{3}\right) \rightarrow \mathfrak{X}\left(M^{3}\right)$ defined by

$$
P_{k}=\sum_{j=0}^{k}(-1)^{j}\binom{3}{k-j} H_{k-j} S^{j}=(-1)^{k} \sum_{j=0}^{k} a_{k-j} S^{j} .
$$

In particular,
(2) $\quad P_{0}=I, \quad P_{1}=3 H I-S, \quad P_{2}=3 H_{2} I-S \circ P_{1}, \quad P_{3}=H_{3} I-S \circ P_{2}$.

Note that by Cayley-Hamilton theorem we have $P_{3}=0$. Let us recall that, for every point $p \in M^{3}$, each $P_{k}(p)$ is also a self-adjoint linear operator on the tangent hyperplane $T_{p} M$ which commutes with $S(p)$. Indeed, $S(p)$ and $P_{k}(p)$ can be simultaneously diagonalized: if $\left\{e_{1}, e_{2}, e_{3}\right\}$ are the eigenvectors of $S(p)$ corresponding to the eigenvalues $\kappa_{1}(p), \kappa_{2}(p), \kappa_{3}(p)$, respectively, then they are also the eigenvectors of $P_{k}(p)$ with corresponding eigenvalues given by

$$
\mu_{k}^{i}(p)=\sum_{\substack{i_{1}<\cdots<i_{1} \\ i_{j} \notin i}}^{3} \kappa_{i_{1}} \cdots \kappa_{i_{k}}, \quad \text { for every } i=1,2,3 \text { and } k=1,2 .
$$

In particular,

$$
\begin{array}{lll}
\mu_{1}^{1}=\kappa_{2}+\kappa_{3}, & \mu_{1}^{2}=\kappa_{1}+\kappa_{3}, & \mu_{1}^{3}=\kappa_{1}+\kappa_{2} \\
\mu_{2}^{1}=\kappa_{2} \kappa_{3}, & \mu_{2}^{2}=\kappa_{1} \kappa_{3}, & \mu_{2}^{3}=\kappa_{1} \kappa_{2}
\end{array}
$$

According to [22, p. 86], the divergence of a vector field $X$ is the differentiable function defined as the contraction of the operator $\nabla X$, where $\nabla X(Y):=\nabla_{Y} X$, that is,

$$
\operatorname{div}(X)=C(\nabla X)=\operatorname{tr}(\nabla X)=\sum_{i, j} g^{i j}\left\langle\nabla_{E_{i}} X, E_{j}\right\rangle
$$

$\left\{E_{i}\right\}$ being any local frame of tangent vectors fields, where $\left(g^{i j}\right)$ represents the inverse of the metric $\left(g_{i j}\right)=\left(\left\langle E_{i}, E_{j}\right\rangle\right)$. For an operator $T: \mathfrak{X}\left(M^{3}\right) \longrightarrow \mathfrak{X}\left(M^{3}\right)$ we have two divergences: one associated to the $(1,1)$-contraction $C_{1}^{1}$, and another associated to the metric contraction $C_{12}$; the first contraction produces a 1-form and the second contraction produces a vector field. We consider here the second one, so that the divergence of an operator $T$ will be the vector field $\operatorname{div}(T) \in \mathfrak{X}\left(M^{3}\right)$ defined as

$$
\operatorname{div}(T)=C_{12}(\nabla T)=\sum_{i, j} g^{i j}\left(\nabla_{E_{i}} T\right) E_{j}
$$

where $\nabla T(X, Y)=\left(\nabla_{X} T\right) Y=\nabla_{X}(T Y)-T\left(\nabla_{X} Y\right)$.
In the following lemma (see [19] for details) we present some interesting properties of the Newton transformations. The proof of the first four is merely algebraic and straightforward.

Lemma 1. The Newton transformations $P_{k}, k=1,2$, satisfy the following properties:
(a) $\operatorname{tr}\left(P_{k}\right)=c_{k} H_{k}$,
(b) $\operatorname{tr}\left(S \circ P_{k}\right)=c_{k} H_{k+1}$,
(c) $\operatorname{tr}\left(S^{2} \circ P_{1}\right)=9 H H_{2}-3 H_{3}$,
(d) $\operatorname{tr}\left(S^{2} \circ P_{2}\right)=3 H H_{3}$,
(e) $\operatorname{tr}\left(\nabla_{X} S \circ P_{k}\right)=\binom{3}{k+1}\left\langle\nabla H_{k+1}, X\right\rangle$,
$(f) \operatorname{div}\left(P_{k}\right)=0$,
where $c_{1}=6$ and $c_{2}=3$.
Keeping in mind this lemma we obtain

$$
\operatorname{div}\left(P_{k}(\nabla f)\right)=\operatorname{tr}\left(P_{k} \circ \nabla^{2} f\right)
$$

where $\nabla^{2} f: \mathfrak{X}\left(M^{3}\right) \longrightarrow \mathfrak{X}\left(M^{3}\right)$ denotes the self-adjoint linear operator metrically equivalent to the Hessian of $f$, given by $\left\langle\nabla^{2} f(X), Y\right\rangle=\left\langle\nabla_{X}(\nabla f), Y\right\rangle$, for vector
fields $X, Y \in \mathfrak{X}\left(M^{3}\right)$. Associated to each Newton transformation $P_{k}$, we can define the second-order linear differential operator $L_{k}: \mathcal{C}^{\infty}\left(M^{3}\right) \longrightarrow \mathcal{C}^{\infty}\left(M^{3}\right)$ by

$$
\begin{equation*}
L_{k}(f)=\operatorname{tr}\left(P_{k} \circ \nabla^{2} f\right) \tag{3}
\end{equation*}
$$

An interesting property of $L_{k}$ is the following:

$$
\begin{equation*}
L_{k}(f g)=g L_{k}(f)+f L_{k}(g)+2\left\langle P_{k}(\nabla f), \nabla g\right\rangle \tag{4}
\end{equation*}
$$

for every couple of differentiable functions $f, g \in C^{\infty}\left(M^{3}\right)$.

## 3. First Formulas and Examples

First we will calculate $L_{k}$ acting on the coordinate components of the immersion $\psi$, that is, a function given by $\langle\psi, e\rangle$, where $e \in \mathbb{R}_{1}^{5}$ is an arbitrary fixed vector. An easy computation shows that

$$
\begin{equation*}
\nabla\langle\psi, e\rangle=e^{\top}=e-\langle N, e\rangle N+\langle\psi, e\rangle \psi \tag{5}
\end{equation*}
$$

where $e^{\top} \in \mathfrak{X}\left(M^{3}\right)$ denotes the tangential component of $e$. Taking covariant derivative in (5), and using the Gauss and Weingarten formulae, we obtain

$$
\begin{equation*}
\nabla_{X} \nabla\langle\psi, e\rangle=\nabla_{X} e^{\top}=\langle N, e\rangle S X+\langle\psi, e\rangle X \tag{6}
\end{equation*}
$$

for every vector field $X \in \mathfrak{X}\left(M^{3}\right)$. Finally, by using (3) and Lemma 1, we obtain

$$
\begin{equation*}
L_{k}\langle\psi, e\rangle=c_{k} H_{k+1}\langle N, e\rangle+c_{k} H_{k}\langle\psi, e\rangle . \tag{7}
\end{equation*}
$$

This expression allows us to extend operator $L_{k}$ to vector functions $F=\left(f_{1}, \ldots, f_{5}\right)$, $f_{i} \in \mathcal{C}^{\infty}\left(M^{3}\right)$, as follows: $L_{k} F:=\left(L_{k} f_{1}, \ldots, L_{k} f_{5}\right)$. Then $L_{k} \psi$ can be computed as

$$
\begin{equation*}
L_{k} \psi=c_{k} H_{k+1} N+c_{k} H_{k} \psi \tag{8}
\end{equation*}
$$

where $\left\{e_{1}, \ldots, e_{5}\right\}$ stands for an orthonormal basis in $\mathbb{R}_{1}^{5}$.
Now, we will compute $L_{k} N$, and in order to do that we are going to compute the operator $L_{k}$ acting on the coordinate functions of the Gauss map $N$, that is, the functions $\langle N, e\rangle$ where $e \in \mathbb{R}_{1}^{5}$ is an arbitrary fixed vector. A straightforward computation yields

$$
\nabla\langle N, e\rangle=-S e^{\top}
$$

that jointly with the Weingarten formula and (6), leads to

$$
\nabla_{X} \nabla\langle N, e\rangle=-\left(\nabla_{e^{\top}} S\right) X-\langle N, e\rangle S^{2} X-\langle\psi, e\rangle S X
$$

for every tangent vector field $X$. This equation, combined with (3) and Lemma 1, yields

$$
\begin{align*}
L_{k}\langle N, e\rangle & =-\operatorname{tr}\left(\nabla_{e^{\top}} S \circ P_{k}\right)-\langle N, e\rangle \operatorname{tr}\left(S^{2} \circ P_{k}\right)-\langle\psi, e\rangle \operatorname{tr}\left(S \circ P_{k}\right) \\
& =-\binom{3}{k+1}\left\langle\nabla H_{k+1}, e\right\rangle-\operatorname{tr}\left(S^{2} \circ P_{k}\right)\langle N, e\rangle-c_{k} H_{k+1}\langle\psi, e\rangle, \tag{9}
\end{align*}
$$

which is equivalent to

$$
L_{k} N=-\binom{3}{k+1} \nabla H_{k+1}-\operatorname{tr}\left(S^{2} \circ P_{k}\right) N-c_{k} H_{k+1} \psi
$$

On the other hand, equations (4) and (7) lead to

$$
\begin{aligned}
L_{k}^{2}\langle\psi, e\rangle= & c_{k} H_{k+1} L_{k}\langle N, e\rangle+L_{k}\left(c_{k} H_{k+1}\right)\langle N, e\rangle+2 c_{k}\left\langle P_{k}\left(\nabla H_{k+1}\right), \nabla\langle N, e\rangle\right\rangle \\
& +c_{k} H_{k} L_{k}\langle\psi, e\rangle+L_{k}\left(c_{k} H_{k}\right)\langle\psi, e\rangle+2 c_{k}\left\langle P_{k}\left(\nabla H_{k}\right), \nabla\langle\psi, e\rangle\right\rangle,
\end{aligned}
$$

and by using again (7) and (9) we get

$$
\begin{aligned}
& L_{k}^{2}\langle\psi, e\rangle \\
= & -c_{k}\binom{3}{k+1} H_{k+1}\left\langle\nabla H_{k+1}, e\right\rangle-2 c_{k}\left\langle\left(S \circ P_{k}\right)\left(\nabla H_{k+1}\right), e\right\rangle+2 c_{k}\left\langle P_{k}\left(\nabla H_{k}\right), e\right\rangle \\
& +\left[c_{k} L_{k}\left(H_{k+1}\right)-\left(\operatorname{tr}\left(S^{2} \circ P_{k}\right)-c_{k} H_{k}\right) c_{k} H_{k+1}\right]\langle N, e\rangle \\
& +\left[-c_{k}^{2} H_{k+1}^{2}+c_{k}^{2} H_{k}^{2}+c_{k} L_{k}\left(H_{k}\right)\right]\langle\psi, e\rangle .
\end{aligned}
$$

Finally, we obtain

$$
\begin{align*}
L_{k}^{2} \psi= & -\frac{c_{k}}{2}\binom{3}{k+1} \nabla H_{k+1}^{2}-2 c_{k}\left(S \circ P_{k}\right)\left(\nabla H_{k+1}\right)+2 c_{k} P_{k}\left(\nabla H_{k}\right) \\
& +\left[c_{k} L_{k}\left(H_{k+1}\right)-\left(\operatorname{tr}\left(S^{2} \circ P_{k}\right)-c_{k} H_{k}\right) c_{k} H_{k+1}\right] N  \tag{10}\\
& +\left[-c_{k}^{2} H_{k+1}^{2}+c_{k}^{2} H_{k}^{2}+c_{k} L_{k}\left(H_{k}\right)\right] \psi .
\end{align*}
$$

Example 1. $k$-minimal $H_{k}$-hypersurfaces in $\mathbb{H}^{4}$ are of $L_{k}$-1-type or $L_{k}$-null-1type. In fact, from (8) we obtain that $L_{k} \psi=\lambda \psi$, with $\lambda=c_{k} H_{k}$, and then $M^{3}$ is of $L_{k}$-1-type if $H_{k} \neq 0$; otherwise, $M^{3}$ is of $L_{k}$-null-1-type.

Example 2. Non-flat totally umbilical hypersurfaces in $\mathbb{H}^{4}$ are of $L_{k}$-1-type. As is well known, totally umbilical hypersurfaces in $\mathbb{H}^{4}$ are obtained as the intersection of $\mathbb{H}^{4}$ with a hyperplane of $\mathbb{R}_{1}^{5}$, and the causal character of the hyperplane determines the type of the hypersurface. More precisely, let $a \in \mathbb{R}_{1}^{5}$ be a non-zero constant vector with $\langle a, a\rangle \in\{1,0,-1\}$, and take the differentiable function $f_{a}: \mathbb{H}^{4} \subset \mathbb{R}_{1}^{5} \rightarrow \mathbb{R}$ defined by
$f_{a}(x)=\langle x, a\rangle$. It is not difficult to see that for every $\tau \in \mathbb{R}$ with $\langle a, a\rangle+\tau^{2}=\delta^{2}>0$, the set

$$
M_{\tau}=f_{a}^{-1}(\tau)=\left\{x \in \mathbb{H}^{4} \mid\langle x, a\rangle=\tau\right\}
$$

is a totally umbilical hypersurface in $\mathbb{H}^{4}$, with Gauss map

$$
N(x)=\frac{1}{\delta}(a+\tau x),
$$

and shape operator

$$
\begin{equation*}
S X=-\frac{\tau}{\delta} X \tag{11}
\end{equation*}
$$

It is easy to see, from (11), that $M_{\tau}$ has constant mean curvature $H=-\tau / \delta$ and constant Gauss-Kronecker curvature $K=-1+H^{2}=-\langle a, a\rangle \delta^{-2}$. Therefore, $H_{k}$ and $H_{k+1}$ are also nonzero constants.

Now we will consider all different possibilities:
(i) If $\langle a, a\rangle=-1$, then $|\tau|>1, K=1 /\left(\tau^{2}-1\right)$ is positive, and $M_{\tau} \equiv$ $\mathbb{S}^{3}\left(\sqrt{\tau^{2}-1}\right)$.
(ii) If $\langle a, a\rangle=0$, then $\tau \neq 0, K=0$, and $M_{\tau} \equiv \mathbb{R}^{3}$.
(iii) If $\langle a, a\rangle=1$, then $K=-1 /\left(\tau^{2}+1\right)$ is negative, and $M_{\tau} \equiv \mathbb{H}^{3}\left(-\sqrt{\tau^{2}+1}\right)$.

Bearing (8) in mind we find that $L_{k} \psi=\lambda \psi+b$, where $\lambda=c_{k} H^{k}\left(1-H^{2}\right)$ and $b=c_{k} H^{k+1} \delta^{-1} a$. We distinguish three cases:
(i) If $H=0$, then $M^{3}$ is of $L_{k}$-null-1-type.
(ii) If $|H|=1$, then $\langle a, a\rangle=0$ and $M^{3}$ is flat.
(iii) Otherwise, $\lambda \neq 0$ and we can write

$$
\psi=\psi_{0}+\psi_{1}, \quad \psi_{0}=-\frac{b}{\lambda} \quad \text { and } \quad \psi_{1}=\psi+\frac{b}{\lambda},
$$

where $\psi_{0}$ is constant and $L_{k} \psi_{1}=\lambda \psi_{1}$. Therefore, $M^{3}$ is $L_{k}$-1-type in $\mathbb{R}_{1}^{5}$.
The following proposition shows that the hypersurfaces exhibited in Examples 1 and 2 are the only hypersurfaces in $\mathbb{H}^{4}$ of $L_{k}$-1-type in $\mathbb{R}_{1}^{5}$.

Proposition 2. $k$-minimal $H_{k}$-hypersurfaces in $\mathbb{H}^{4}$ and open portions of a non-flat totally umbilical hypersurface in $\mathbb{H}^{4}$ are the only $L_{k}$-1-type hypersurfaces in $\mathbb{H}^{4}$.

Proof. Let $M^{3}$ be a $L_{k}$-1-type hypersurface in $\mathbb{H}^{4}$, then its position vector $\psi$ can be put as $\psi=\psi_{0}+\psi_{1}$, where $\psi_{0}$ is a constant vector and $L_{k} \psi_{1}=\lambda \psi_{1}$. Hence we deduce $L_{k} \psi=\lambda \psi+b$, with $b=-\lambda \psi_{0}$. From (8) we get

$$
b=c_{k} H_{k+1} N+\left(c_{k} H_{k}-\lambda\right) \psi,
$$

and taking covariant derivative here we obtain

$$
0=-c_{k} H_{k+1} S X+\left(c_{k} H_{k}-\lambda\right) X+c_{k} X\left(H_{k+1}\right) N+c_{k} X\left(H_{k}\right) \psi,
$$

for every vector field $X \in \mathfrak{X}\left(M^{3}\right)$. The previous equation implies that $H_{k}$ and $H_{k+1}$ are both constant. If $H_{k+1} \neq 0$ then we get $S X=\mu X$, for a certain constant $\mu$, i.e. $M^{3}$ is totally umbilical, and then the result follows from Example 2.

Example 3. Standard Riemannian products $\mathbb{H}^{1}\left(r_{1}\right) \times \mathbb{S}^{2}\left(r_{2}\right)$ and $\mathbb{H}^{2}\left(r_{1}\right) \times \mathbb{S}^{1}\left(r_{2}\right)$, with $-r_{1}^{2}+r_{2}^{2}=-1$, are hypersurfaces in $\mathbb{H}^{4}$ of $L_{k}$-2-type in $\mathbb{R}_{1}^{5}$.

For a positive number $r$, let us denote $M_{m}^{3}(r)=\mathbb{H}^{m}\left(-\sqrt{1+r^{2}}\right) \times \mathbb{S}^{3-m}(r) \subset \mathbb{H}^{4}$, $m=1,2$. In the case $m=1$, observe that the hypersurface $M_{1}^{3}(r)$ is defined by the equation

$$
M_{1}^{3}(r)=\left\{x \in \mathbb{H}^{4} \mid x_{3}^{2}+x_{4}^{2}+x_{5}^{2}=r^{2}\right\},
$$

and its Gauss map is given by

$$
N(x)=\left(\frac{r}{\sqrt{1+r^{2}}} x_{1}, \frac{r}{\sqrt{1+r^{2}}} x_{2}, \frac{\sqrt{1+r^{2}}}{r} x_{3}, \frac{\sqrt{1+r^{2}}}{r} x_{4}, \frac{\sqrt{1+r^{2}}}{r} x_{5}\right) .
$$

Then its principal curvatures in $\mathbb{H}^{4}$ are

$$
\kappa_{1}=\frac{-r}{\sqrt{1+r^{2}}} \quad \text { and } \quad \kappa_{2}=\kappa_{3}=\frac{-\sqrt{1+r^{2}}}{r} .
$$

Hence we get

$$
H_{1}=-\frac{2+3 r^{2}}{3 r \sqrt{1+r^{2}}}, \quad H_{2}=\frac{1+3 r^{2}}{3 r^{2}}, \quad H_{3}=-\frac{\sqrt{1+r^{2}}}{r} .
$$

If we put $\psi_{1}=\left(x_{1}, x_{2}, 0,0,0\right)$ and $\psi_{2}=\left(0,0, x_{3}, x_{4}, x_{5}\right)$, then $\psi=\psi_{1}+\psi_{2}$ and by using (8) we obtain:
(a) $L_{0} \psi_{1}=\lambda_{1} \psi_{1}$ and $L_{0} \psi_{2}=\lambda_{2} \psi_{2}$, where $\lambda_{1}=\frac{1}{1+r^{2}}$ and $\lambda_{2}=-\frac{2}{r^{2}}$. Therefore, $M_{1}^{3}(r)$ is of $L_{0}-2$-type in $\mathbb{R}_{1}^{5}$ for any $r$ (see [11, Example 1]).
(b) $L_{1} \psi_{1}=\lambda_{1} \psi_{1}$ and $L_{1} \psi_{2}=\lambda_{2} \psi_{2}$, where $\lambda_{1}=-\frac{2}{r \sqrt{1+r^{2}}}$ and $\lambda_{2}=\frac{2\left(1+2 r^{2}\right)}{r^{3} \sqrt{1+r^{2}}}$. Therefore, $M_{1}^{3}(r)$ is of $L_{1}$-2-type in $\mathbb{R}_{1}^{5}$ for any $r$.
(c) $L_{2} \psi_{1}=\lambda_{1} \psi_{1}$ and $L_{2} \psi_{2}=\lambda_{2} \psi_{2}$, where $\lambda_{1}=\frac{1}{r^{2}}$ and $\lambda_{2}=-\frac{2}{r^{2}}$. Therefore, $M_{1}^{3}(r)$ is of $L_{2}$-2-type in $\mathbb{R}_{1}^{5}$ for any $r$.

In the case $m=2$, note that the hypersurface $M_{2}^{3}(r)$ is defined by the equation

$$
M_{2}^{3}(r)=\left\{x \in \mathbb{H}^{4} \mid x_{4}^{2}+x_{5}^{2}=r^{2}\right\} .
$$

In this case, the Gauss map on $M_{2}^{3}(r)$ in $\mathbb{H}^{4}$ is given by

$$
N(x)=\left(\frac{r}{\sqrt{1+r^{2}}} x_{1}, \frac{r}{\sqrt{1+r^{2}}} x_{2}, \frac{r}{\sqrt{1+r^{2}}} x_{3}, \frac{\sqrt{1+r^{2}}}{r} x_{4}, \frac{\sqrt{1+r^{2}}}{r} x_{5}\right),
$$

and its principal curvatures in $\mathbb{H}^{4}$ are

$$
\kappa_{1}=\kappa_{2}=\frac{-r}{\sqrt{1+r^{2}}} \quad \text { and } \quad \kappa_{3}=\frac{-\sqrt{1+r^{2}}}{r}
$$

Consequently, we get

$$
H_{1}=-\frac{1+3 r^{2}}{3 r \sqrt{1+r^{2}}}, \quad H_{2}=\frac{2+3 r^{2}}{3\left(1+r^{2}\right)}, \quad H_{3}=-\frac{r}{\sqrt{1+r^{2}}} .
$$

If we put as before $\psi_{1}=\left(x_{1}, x_{2}, x_{3}, 0,0\right)$ and $\psi_{2}=\left(0,0,0, x_{4}, x_{5}\right)$, then $\psi=\psi_{1}+\psi_{2}$ and by using (8) we obtain:
(a) $L_{0} \psi_{1}=\lambda_{1} \psi_{1}$ and $L_{0} \psi_{2}=\lambda_{2} \psi_{2}$, where $\lambda_{1}=\frac{2}{1+r^{2}}$ and $\lambda_{2}=-\frac{1}{r^{2}}$. Therefore, $M_{2}^{3}(r)$ is of $L_{0}-2$-type in $\mathbb{R}_{1}^{5}$ for any $r$ (see [11, Example 1]).
(b) $L_{1} \psi_{1}=\lambda_{1} \psi_{1}$ and $L_{1} \psi_{2}=\lambda_{2} \psi_{2}$, where $\lambda_{1}=-\frac{2\left(1+2 r^{2}\right)}{r\left(1+r^{2}\right)^{3 / 2}}$ and $\lambda_{2}=\frac{2}{r \sqrt{1+r^{2}}}$. Therefore, $M_{2}^{3}(r)$ is of $L_{1}$-2-type in $\mathbb{R}_{1}^{5}$ for any $r$.
(c) $L_{2} \psi_{1}=\lambda_{1} \psi_{1}$ and $L_{2} \psi_{2}=\lambda_{2} \psi_{2}$, where $\lambda_{1}=\frac{2}{1+r^{2}}$ and $\lambda_{2}=-\frac{1}{1+r^{2}}$. Therefore, $M_{2}^{3}(r)$ is of $L_{2}$-2-type in $\mathbb{R}_{1}^{5}$ for any $r$.

## 4. The Three-dimensional Case

Let us suppose that a hypersurface $M^{3}$ in $\mathbb{H}^{4}$ is of $L_{k}$-2-type in $\mathbb{R}_{1}^{5}$, that is, its position vector $\psi$ can be written as follows

$$
\psi=a+\psi_{1}+\psi_{2}, \quad L_{k} \psi_{1}=\lambda_{1} \psi_{1}, \quad L_{k} \psi_{2}=\lambda_{2} \psi_{2},
$$

where $a$ is a constant vector in $\mathbb{R}_{1}^{5}$ and $\psi_{1}, \psi_{2}$ are $\mathbb{R}_{1}^{5}$-valued non-constant differentiable functions defined on $M^{3}$.

It is easy to see that $L_{k} \psi=\lambda_{1} \psi_{1}+\lambda_{2} \psi_{2}$ and $L_{k}^{2} \psi=\lambda_{1}^{2} \psi_{1}+\lambda_{2}^{2} \psi_{2}$, and thus

$$
L_{k}^{2} \psi=\left(\lambda_{1}+\lambda_{2}\right) L_{k} \psi-\lambda_{1} \lambda_{2}(\psi-a) .
$$

By using (8) we get

$$
\begin{aligned}
L_{k}^{2} \psi= & \lambda_{1} \lambda_{2} a^{\top}+\left[\left(\lambda_{1}+\lambda_{2}\right) c_{k} H_{k+1}+\lambda_{1} \lambda_{2}\langle N, a\rangle\right] N \\
& +\left[\left(\lambda_{1}+\lambda_{2}\right) c_{k} H_{k}-\lambda_{1} \lambda_{2}-\lambda_{1} \lambda_{2}\langle\psi, a\rangle\right] \psi,
\end{aligned}
$$

that, jointly with (10), yields the following equations of $L_{k}$ - 2 -type,

$$
\begin{align*}
\lambda_{1} \lambda_{2} a^{\top} & =-\frac{c_{k}}{2}\binom{3}{k+1} \nabla H_{k+1}^{2}-2 c_{k}\left(S \circ P_{k}\right)\left(\nabla H_{k+1}\right)+2 c_{k} P_{k}\left(\nabla H_{k}\right),  \tag{12}\\
\lambda_{1} \lambda_{2}\langle N, a\rangle & =c_{k} L_{k}\left(H_{k+1}\right)-\left(\operatorname{tr}\left(S^{2} \circ P_{k}\right)-c_{k} H_{k}+\lambda_{1}+\lambda_{2}\right) c_{k} H_{k+1},  \tag{13}\\
\lambda_{1} \lambda_{2}\langle\psi, a\rangle & =c_{k}^{2} H_{k+1}^{2}-\left(c_{k} H_{k}-\lambda_{1}\right)\left(c_{k} H_{k}-\lambda_{2}\right)-c_{k} L_{k}\left(H_{k}\right) . \tag{14}
\end{align*}
$$

In [13], the author shows that if $M^{n}$ is a hypersurface of the hyperbolic space $\mathbb{H}^{n+1}$ with constant mean curvature and constant scalar curvature, then $M^{n}$ is either of 1-type or of 2-type. He also proves that every 2-type hypersurface of the hyperbolic space has nonzero constant mean curvature and constant scalar curvature.

Our goal in this section is to prove similar results for operators $L_{1}$ and $L_{2}$.
Theorem 3. Let $\psi: M^{3} \rightarrow \mathbb{H}^{4} \subset \mathbb{R}_{1}^{5}$ be an orientable $H_{2}$-hypersurface. If $M^{3}$ is of $L_{2}$-2-type then the Gauss-Kronecker curvature $H_{3}$ is a nonzero constant.

Proof. Let $\left\{E_{1}, E_{2}, E_{3}\right\}$ be a local orthonormal frame of principal directions of $S$ such that $S E_{i}=\kappa_{i} E_{i}$ for every $i=1,2,3$, and consider the open set

$$
\mathcal{U}_{3}=\left\{p \in M^{3} \mid \nabla H_{3}^{2}(p) \neq 0\right\} .
$$

Let us suppose that $\mathcal{U}_{3}$ is not empty. Since we are assuming that $M^{3}$ is of $L_{2}-2$-type and $H_{2}$ is constant, then by taking covariant derivative in (14) we have $\lambda_{1} \lambda_{2} a^{\top}=9 \nabla H_{3}^{2}$, and putting this into (12) yields

$$
\begin{equation*}
\left(S \circ P_{2}\right)\left(\nabla H_{3}^{2}\right)=-\frac{7}{2} H_{3} \nabla H_{3}^{2} \quad \text { on } \mathcal{U}_{3} . \tag{15}
\end{equation*}
$$

Since $P_{3}=0$ then $S \circ P_{2}=H_{3} I$ and so $\left(S \circ P_{2}\right)\left(\nabla H_{3}^{2}\right)=H_{3} \nabla H_{3}^{2}$, that jointly with (15) implies $H_{3} \nabla H_{3}^{2}=0$ on $\mathcal{U}_{3}$, which is not possible.

We want to extend the previous theorem for the operator $L_{1}$; next theorem is an intermediate step.

Recall that a hypersurface $M^{n}$ immersed in either the Euclidean space $\mathbb{R}^{n+1}$, the sphere $\mathbb{S}^{n+1}$ or the hyperbolic space $\mathbb{H}^{n+1}$, is called isoparametric if all the principal curvatures $\kappa_{i}$ are constant functions; this is equivalent to saying that all the mean curvatures $H_{i}$ are constant functions. An isoparametric hypersurface of the Euclidean space can have at most two distinct principal curvatures, and it must be an open portion of a hyperplane, hypersphere or spherical cylinder $\mathbb{S}^{k}(r) \times \mathbb{R}^{n-k}$ (see e.g. [26, 25]). A similar result holds for $\mathbb{H}^{n+1}$ : an isoparametric hypersurface must be an open part of a totally umbilical hypersurface or hyperbolic cylinder $\mathbb{H}^{m}\left(r_{1}\right) \times \mathbb{S}^{n-m}\left(r_{2}\right)$ (see [3]). However, the classification of isoparametric hipersurfaces in the sphere $\mathbb{S}^{n+1}$ turns out to be much more complicated, as Elie Cartan showed (see [4, 5, 6]).

Theorem 4. Let $M^{3}$ be an orientable $H_{k}$-hypersurface of the hyperbolic space $\mathbb{H}^{4}$, which is not totally umbilical, and consider the following three conditions:
(a) $H_{k+1}$ is a nonzero constant.
(b) $\operatorname{tr}\left(S^{2} \circ P_{k}\right)$ is constant.
(c) $M^{3}$ is of $L_{k}$-2-type.

Then any two conditions imply the third one.
Proof. First, we show that conditions $a$ ) and $b$ ) imply condition $c$ ). From Lemma 1 we obtain that $M^{3}$ is an isoparametric hypersurface; since $M^{3}$ is not totally umbilical then $M^{3}$ is a hyperbolic cylinder, and then the claim follows from Example 3.

Secondly, we show that conditions $a$ ) and $c$ ) imply condition $b$ ). By taking covariant differentiation in equation (13), and bearing (14) in mind, we find

$$
c_{k} H_{k+1} X\left(\operatorname{tr}\left(S^{2} \circ P_{k}\right)\right)=-\lambda_{1} \lambda_{2} X(\langle N, a\rangle)=\lambda_{1} \lambda_{2}\left\langle S X, a^{\top}\right\rangle=0,
$$

that is, $\operatorname{tr}\left(S^{2} \circ P_{k}\right)$ is constant on $M^{3}$.
Finally, we show that conditions $b$ ) and $c$ ) imply condition $a$ ). In the case $k=2$, the proof follows directly from Theorem 3. To prove the claim in the case $k=1$, let us consider the open set

$$
\mathcal{U}_{2}=\left\{p \in M^{3} \mid \nabla H_{2}^{2}(p) \neq 0\right\},
$$

and assume that it is not empty. Since $H$ is constant, by taking covariant derivative in (14) we obtain that $\lambda_{1} \lambda_{2} a^{\top}=36 \nabla H_{2}^{2}$. Using this in (12) we get

$$
\begin{equation*}
\left(S \circ P_{1}\right)\left(\nabla H_{2}^{2}\right)=-\frac{15}{2} H_{2} \nabla H_{2}^{2} \quad \text { on } \mathcal{U}_{2}, \tag{16}
\end{equation*}
$$

that jointly with equation (2) leads to $P_{2}\left(\nabla H_{2}^{2}\right)=\frac{21}{2} H_{2} \nabla H_{2}^{2}$. Now, by applying the operator $S$ on both sides, we have

$$
\begin{equation*}
\left(S \circ P_{2}\right)\left(\nabla H_{2}^{2}\right)=\frac{21}{2} H_{2} S\left(\nabla H_{2}^{2}\right) . \tag{17}
\end{equation*}
$$

Since $P_{3}=0$ we get $S \circ P_{2}=H_{3} I$, and then $\left(S \circ P_{2}\right)\left(\nabla H_{2}^{2}\right)=H_{3} \nabla H_{2}^{2}$, that jointly with (17) implies

$$
S\left(\nabla H_{2}^{2}\right)=\frac{2 H_{3}}{21 H_{2}} \nabla H_{2}^{2}
$$

Without loss of generality, let us assume that $E_{1}$ is parallel to $\nabla H_{2}^{2}$, i.e. the principal curvature $\kappa_{1}=\frac{2 H_{3}}{21 H_{2}}$. Then we have

$$
\left(S \circ P_{1}\right)\left(\nabla H_{2}^{2}\right)=\kappa_{1} \mu_{1}^{1} \nabla H_{2}^{2}=\frac{2 H_{3}}{21 H_{2}}\left(3 H-\frac{2 H_{3}}{21 H_{2}}\right) \nabla H_{2}^{2},
$$

that jointly with (16) yields the following equation,

$$
6615 H_{2}^{3}+252 H_{2} H_{3}-8 H_{3}^{2}=0
$$

From Lemma 1 we have that $3 H_{3}=9 H H_{2}-\operatorname{tr}\left(S^{2} \circ P_{1}\right)$, and then the previous equation can be rewritten as follows

$$
6615 H_{2}^{3}+684 H^{2} H_{2}^{2}-68 H \operatorname{tr}\left(S^{2} \circ P_{1}\right) H_{2}-\frac{8}{9} \operatorname{tr}\left(S^{2} \circ P_{1}\right)=0
$$

In other words, $\mathrm{H}_{2}$ is a root of a polynomial with constant coefficients, and so $\mathrm{H}_{2}$ has to be constant, which is a contradiction.

An interesting consequence of the last theorem is the following result.
Theorem 5. Let $\psi: M^{3} \rightarrow \mathbb{H}^{4} \subset \mathbb{R}_{1}^{5}$ be an orientable $H_{2}$-hypersurface. If $M$ is of $L_{2}$-2-type then $M^{3}$ is an isoparametric hypersurface.

Proof. From Theorem 3 we get that $H_{3}$ is a nonzero constant, and then Theorem 4 yields that $\operatorname{tr}\left(S^{2} \circ P_{2}\right)$ is constant. Now we use Lemma $1(\mathrm{~d})$ to deduce that the mean curvature $H$ is constant, and this concludes the proof.

Since the isoparametric hypersurfaces of the hyperbolic space $\mathbb{H}^{4} \subset \mathbb{R}_{1}^{5}$ are well known, the following result is clear.

Theorem 6. Let $\psi: M^{3} \rightarrow \mathbb{H}^{4} \subset \mathbb{R}_{1}^{5}$ be an orientable $H_{2}$-hypersurface, which is not totally umbilical. Then $M^{3}$ is of $L_{2}$-2-type if and only if $M^{3}$ is a standard Riemannian product $\mathbb{H}^{1}\left(r_{1}\right) \times \mathbb{S}^{2}\left(r_{2}\right)$ or $\mathbb{H}^{2}\left(r_{1}\right) \times \mathbb{S}^{1}\left(r_{2}\right)$, with $-r_{1}^{2}+r_{2}^{2}=-1$.

Now, we state the main result of this section.
Theorem 7. Let $\psi: M^{3} \rightarrow \mathbb{H}^{4} \subset \mathbb{R}_{1}^{5}$ be an orientable $H_{k}$-hypersurface. If $M$ is of $L_{k}$-2-type then $H_{k+1}$ is a nonzero constant.

Proof. Case $k=0$ is shown in [13] and case $k=2$ has been proved in Theorem 3, so we can assume $k=1$. Let us consider $\left\{E_{1}, E_{2}, E_{3}\right\}$ a local orthonormal frame of principal directions of $S$ such that $S E_{i}=\kappa_{i} E_{i}$ for every $i=1,2,3$. Let us define the open set

$$
\mathcal{U}_{2}=\left\{p \in M^{3} \mid \nabla H_{2}^{2}(p) \neq 0\right\}
$$

and suppose that $\mathcal{U}_{2}$ is not empty. Since we are assuming that $M^{3}$ is $L_{1}$-2-type and $H$ is constant, then equation (14) leads to

$$
\begin{equation*}
\lambda_{1} \lambda_{2} a^{\top}=36 \nabla H_{2}^{2} \tag{18}
\end{equation*}
$$

Using this equation in (12) we have that $\left(S \circ P_{1}\right)\left(\nabla H_{2}^{2}\right)=-\frac{15}{2} H_{2} \nabla H_{2}^{2}$ on $\mathcal{U}_{2}$, and substituting this into (2) we obtain

$$
\begin{equation*}
P_{2}\left(\nabla H_{2}^{2}\right)=\frac{21}{2} H_{2} \nabla H_{2}^{2} \quad \text { on } \mathcal{U}_{2} . \tag{19}
\end{equation*}
$$

The vector field $\nabla H_{2}^{2}$ can be written as $\nabla H_{2}^{2}=E_{1}\left(H_{2}^{2}\right) E_{1}+E_{2}\left(H_{2}^{2}\right) E_{2}+E_{3}\left(H_{2}^{2}\right) E_{3}$, and then

$$
P_{2}\left(\nabla H_{2}^{2}\right)=E_{1}\left(H_{2}^{2}\right) \mu_{2}^{1} E_{1}+E_{2}\left(H_{2}^{2}\right) \mu_{2}^{2} E_{2}+E_{3}\left(H_{2}^{2}\right) \mu_{2}^{3} E_{3}
$$

Therefore equation (19) is equivalent to

$$
\begin{equation*}
E_{i}\left(H_{2}^{2}\right)\left(\mu_{2}^{i}-\frac{21}{2} H_{2}\right)=0 \quad \text { on } \mathcal{U}_{2} \tag{20}
\end{equation*}
$$

for every $i=1,2,3$. An immediate and important consequence of this equation is that $E_{i}\left(H_{2}^{2}\right)=0$ for some $i$. Otherwise, we deduce that

$$
\operatorname{tr}\left(P_{2}\right)=\mu_{2}^{1}+\mu_{2}^{2}+\mu_{2}^{3}=\frac{63}{2} H_{2},
$$

that jointly with Lemma 1 leads to $H_{2}=0$ on $\mathcal{U}_{2}$, which is a contradiction.
Bearing in mind the previous consequence, and without loss of generality, we have to analyze the following two possible cases.

Case 1. $E_{1}\left(H_{2}^{2}\right) \neq 0, E_{2}\left(H_{2}^{2}\right) \neq 0$ and $E_{3}\left(H_{2}^{2}\right)=0$.
From (20) we have $\mu_{2}^{1}=\mu_{2}^{2}=\frac{21}{2} H_{2}$, then $\left(\kappa_{1}-\kappa_{2}\right) \kappa_{3}=0$, and therefore $\kappa_{1}=\kappa_{2}$. Observe that $\kappa_{i} \neq 0$ for all $i$, otherwise $H_{2}=0$. It is easy to see that

$$
\kappa_{2} \kappa_{3}=\mu_{2}^{1}=\frac{21}{2} H_{2}=\frac{7}{2}\left(\kappa_{2}^{2}+2 \kappa_{2} \kappa_{3}\right),
$$

and so $7 \kappa_{2}+12 \kappa_{3}=0$. On the other hand, we know that $3 H=2 \kappa_{2}+\kappa_{3}$ and then we get that the principal curvatures $\kappa_{2}$ and $\kappa_{3}$ are constant. So $H_{2}$ is also constant, which can not be possible.

Case 2. $E_{1}\left(H_{2}^{2}\right) \neq 0, E_{2}\left(H_{2}^{2}\right)=0$ and $E_{3}\left(H_{2}^{2}\right)=0$.
We know that $3 H_{2}=\kappa_{1} \mu_{1}^{1}+\mu_{2}^{1}$ and $\mu_{2}^{1}=\frac{21}{2} H_{2}$ (see (20)), then we have

$$
\begin{equation*}
H_{2}=\frac{2}{15}\left(\kappa_{1}^{2}-3 H \kappa_{1}\right) \quad \text { and } \quad H_{2}^{2}=p\left(\kappa_{1}\right) \tag{21}
\end{equation*}
$$

where $p(x)=\left(\frac{2}{15}\right)^{2}\left(x^{4}-6 H x^{3}+9 H^{2} x^{2}\right)$. Observe that $H \neq 0$; otherwise, $\kappa_{2}+\kappa_{3}=$ $-\kappa_{1}$ and from (21) we get $\kappa_{2} \kappa_{3}=\frac{7}{5} \kappa_{1}^{2}$. Then $\kappa_{2}$ and $\kappa_{3}$ are the roots of the equation $t^{2}+\kappa_{1} t+\frac{7}{5} \kappa_{1}^{2}=0$, but this is not possible since the discriminant of this equation is negative.

We claim that

$$
\begin{align*}
E_{1}\left(H_{2}^{2}\right) & =p^{\prime}\left(\kappa_{1}\right) E_{1}\left(\kappa_{1}\right),  \tag{22}\\
\lambda_{1} \lambda_{2}\langle\psi, a\rangle & =36 p\left(\kappa_{1}\right)+A_{0},  \tag{23}\\
\lambda_{1} \lambda_{2}\langle N, a\rangle & =q\left(\kappa_{1}\right)+B_{0}, \tag{24}
\end{align*}
$$

where $q(x)=-\left(\frac{4}{5}\right)^{2}\left(\frac{4}{5} x^{5}-\frac{9 H}{2} x^{4}+6 H^{2} x^{3}\right)$, and $A_{0}, B_{0}$ are two constants. First, (22) and (23) follow directly from (21) and (14), respectively. On the other hand, bearing (18) in mind we find that

$$
\begin{aligned}
X\left(\lambda_{1} \lambda_{2}\langle N, a\rangle\right) & =-\lambda_{1} \lambda_{2}\left\langle S X, a^{\top}\right\rangle=-36 \kappa_{1}\left\langle X, \nabla H_{2}^{2}\right\rangle \\
& =-36 \kappa_{1} X\left(H_{2}^{2}\right)=X\left(q\left(\kappa_{1}\right)\right),
\end{aligned}
$$

for any tangent vector field $X$, and this implies equation (24).
Now, by taking covariant differentiation in (18) in the direction of an arbitrary tangent vector field $X$, we have

$$
\lambda_{1} \lambda_{2} \nabla_{X} a^{\top}=36 X\left(E_{1}\left(H_{2}^{2}\right)\right) E_{1}+36 E_{1}\left(H_{2}^{2}\right) \nabla_{X} E_{1},
$$

that jointly with (6) yields

$$
\begin{equation*}
36 E_{1}\left(H_{2}^{2}\right) \nabla_{X} E_{1}=-36 X\left(E_{1}\left(H_{2}^{2}\right)\right) E_{1}+\lambda_{1} \lambda_{2}(\langle N, a\rangle S X+\langle\psi, a\rangle X) \tag{25}
\end{equation*}
$$

or equivalently

$$
\begin{align*}
& 36 E_{1}\left(H_{2}^{2}\right)\left\langle\nabla_{X} E_{1}, E_{i}\right\rangle \\
= & -36 X\left(E_{1}\left(H_{2}^{2}\right)\right) \delta_{1 i}+\lambda_{1} \lambda_{2}\left(\langle N, a\rangle \kappa_{i}+\langle\psi, a\rangle\right)\left\langle X, E_{i}\right\rangle, \tag{26}
\end{align*}
$$

for $i=1,2,3$. If we take $X=E_{1}$, then (26) reduces to the following equations

$$
\begin{aligned}
36 E_{1}\left(E_{1}\left(H_{2}^{2}\right)\right) & =\lambda_{1} \lambda_{2}\left(\langle N, a\rangle \kappa_{1}+\langle\psi, a\rangle\right), \\
E_{1}\left(H_{2}^{2}\right)\left\langle\nabla_{E_{1}} E_{1}, E_{i}\right\rangle & =0, \quad i=2,3 .
\end{aligned}
$$

From the last equation we conclude that $\nabla_{E_{1}} E_{1}=0$, that is, the integral curves of $E_{1}$ on $\mathcal{U}_{2}$ are geodesics of $M^{3}$.

Let $X$ be a tangent vector field orthogonal to $E_{1}$. Then equation (26) for $i=1$ leads to $X\left(E_{1}\left(H_{2}^{2}\right)\right)=0$ and thus (25) yields

$$
\begin{equation*}
36 E_{1}\left(H_{2}^{2}\right) \nabla_{X} E_{1}=\lambda_{1} \lambda_{2}(\langle N, a\rangle S X+\langle\psi, a\rangle X), \quad \forall X \perp E_{1} \tag{27}
\end{equation*}
$$

From the Codazzi equation $\left(\nabla_{E_{j}} S\right) E_{1}=\left(\nabla_{E_{1}} S\right) E_{j}$, we get

$$
E_{1}\left(\kappa_{j}\right)=\left(\kappa_{1}-\kappa_{j}\right)\left\langle\nabla_{E_{j}} E_{1}, E_{j}\right\rangle, \quad j=2,3,
$$

that jointly with (27) for $X=E_{j}$ yields

$$
\begin{aligned}
& 36 E_{1}\left(H_{2}^{2}\right) E_{1}\left(\kappa_{j}\right) \\
= & \left(\kappa_{1}-\kappa_{j}\right)\left[\lambda_{1} \lambda_{2}\langle N, a\rangle \kappa_{j}+\lambda_{1} \lambda_{2}\langle\psi, a\rangle\right] \\
= & -\lambda_{1} \lambda_{2}\langle N, a\rangle \kappa_{j}^{2}+\lambda_{1} \lambda_{2}\langle N, a\rangle \kappa_{1} \kappa_{j}-\lambda_{1} \lambda_{2}\langle\psi, a\rangle \kappa_{j}+\lambda_{1} \lambda_{2}\langle\psi, a\rangle \kappa_{1} .
\end{aligned}
$$

Last equation implies

$$
\begin{aligned}
36 E_{1}\left(H_{2}^{2}\right)\left(E_{1}\left(\kappa_{2}\right)+E_{1}\left(\kappa_{3}\right)\right)= & -\lambda_{1} \lambda_{2}\langle N, a\rangle\left(\kappa_{2}^{2}+\kappa_{3}^{2}\right)+\lambda_{1} \lambda_{2}\langle N, a\rangle \kappa_{1}\left(\kappa_{2}+\kappa_{3}\right) \\
& -\lambda_{1} \lambda_{2}\langle\psi, a\rangle\left(\kappa_{2}+\kappa_{3}\right)+2 \lambda_{1} \lambda_{2}\langle\psi, a\rangle \kappa_{1},
\end{aligned}
$$

that is,

$$
\begin{aligned}
36 E_{1}\left(H_{2}^{2}\right) E_{1}\left(3 H-\kappa_{1}\right)= & -\lambda_{1} \lambda_{2}\langle N, a\rangle\left(\operatorname{tr}\left(S^{2}\right)-\kappa_{1}^{2}\right)+\lambda_{1} \lambda_{2}\langle N, a\rangle \kappa_{1}\left(3 H-\kappa_{1}\right) \\
& -\lambda_{1} \lambda_{2}\langle\psi, a\rangle\left(3 H-\kappa_{1}\right)+2 \lambda_{1} \lambda_{2}\langle\psi, a\rangle \kappa_{1} .
\end{aligned}
$$

From (1) and (21) we have that $\operatorname{tr}\left(S^{2}\right)=9 H^{2}-\frac{3}{5} H \kappa_{1}-\frac{4}{5} \kappa_{1}^{2}$. By using this and (22), last equation can be written as

$$
\begin{align*}
& 36 p^{\prime}\left(\kappa_{1}\right)\left[E_{1}\left(\kappa_{1}\right)\right]^{2} \\
= & -\frac{1}{5} \lambda_{1} \lambda_{2}\langle N, a\rangle\left(4 \kappa_{1}^{2}+3 H \kappa_{1}-45 H^{2}\right)+3 \lambda_{1} \lambda_{2}\langle\psi, a\rangle\left(H-\kappa_{1}\right) . \tag{28}
\end{align*}
$$

On the other hand, a direct computation shows

$$
\begin{align*}
36^{2}\left[p^{\prime}\left(\kappa_{1}\right) E_{1}\left(\kappa_{1}\right)\right]^{2} & =36^{2}\left[E_{1}\left(H_{2}^{2}\right)\right]^{2}=36^{2}\left\langle\nabla H_{2}^{2}, \nabla H_{2}^{2}\right\rangle=\lambda_{1}^{2} \lambda_{2}^{2}\left|a^{\top}\right|^{2}  \tag{29}\\
& =\lambda_{1}^{2} \lambda_{2}^{2}|a|^{2}-\left(\lambda_{1} \lambda_{2}\langle N, a\rangle\right)^{2}+\left(\lambda_{1} \lambda_{2}\langle\psi, a\rangle\right)^{2} .
\end{align*}
$$

From equations (28) and (29), and taking into account (23) and (24), we find a polynomial $T(x)$ with constant coefficients given by

$$
\begin{align*}
T(x)= & {\left[q(x)+B_{0}\right]^{2}-\left[36 p(x)+A_{0}\right]^{2} } \\
& -\frac{36}{5}\left[q(x)+B_{0}\right](4 x+15 H)(x-3 H) p^{\prime}(x)  \tag{30}\\
& +108\left[36 p(x)+A_{0}\right](H-x) p^{\prime}(x)-\lambda_{1}^{2} \lambda_{2}^{2}|a|^{2},
\end{align*}
$$

and satisfying $T\left(\kappa_{1}\right)=0$. Therefore, $\kappa_{1}$ is locally constant on $\mathcal{U}_{2}$, and so $H_{2}$ is also constant, which is a contradiction with the definition of $\mathcal{U}_{2}$. This finishes the proof.

An interesting consequence is the following result, similar to Theorem 5.
Theorem 8. Let $\psi: M^{3} \rightarrow \mathbb{H}^{4} \subset \mathbb{R}_{1}^{5}$ be an orientable $H$-hypersurface. If $M^{3}$ is of $L_{1}$-2-type then $M^{3}$ is an isoparametric hypersurface.

Proof. From Theorem 7 we get that $H_{2}$ is a non-zero constant, and then Theorem 4 yields that $\operatorname{tr}\left(S^{2} \circ P_{1}\right)$ is constant. Now we use Lemma 1(c) to deduce that the Gauss-Kronecker curvature $H_{3}$ is constant, and this concludes the proof.

Bearing in mind Theorems 8 and 4, and the classification of isoparametric hypersurfaces in the hyperbolic space $\mathbb{H}^{4}$, the following result, that extends Theorem 6, is clear.

Theorem 9. Let $\psi: M^{3} \rightarrow \mathbb{H}^{4} \subset \mathbb{R}_{1}^{5}$ be an orientable $H$-hypersurface, which is not totally umbilical. Then $M^{3}$ is of $L_{1}$-2-type if and only if $M^{3}$ is a standard Riemannian product $\mathbb{H}^{1}\left(r_{1}\right) \times \mathbb{S}^{2}\left(r_{2}\right)$ or $\mathbb{H}^{2}\left(r_{1}\right) \times \mathbb{S}^{1}\left(r_{2}\right)$, with $-r_{1}^{2}+r_{2}^{2}=-1$.

## 5. The $n$-Dimensional Case

Let $\psi: M^{n} \rightarrow \mathbb{H}^{n+1} \subset \mathbb{R}_{1}^{n+2}$ denote an isometric immersion of an orientable hypersurface $M^{n}$ in the hyperbolic space $\mathbb{H}^{n+1} \equiv \mathbb{H}^{n+1}(0,-1)$. The goal of this section is to classify $L_{k}$-2-type hypersurfaces with constant $k$-th mean curvature $H_{k}$ and having at most two distinct principal curvatures.

Suppose that $\psi$ is of $L_{k}$-2-type, then we can write

$$
\psi=a+\psi_{1}+\psi_{2}, \quad L_{k} \psi_{1}=\lambda_{1} \psi_{1}, \quad L_{k} \psi_{2}=\lambda_{2} \psi_{2}
$$

where $a \in \mathbb{R}_{1}^{n+2}$ is a constant vector and $\psi_{1}, \psi_{2}: M^{n} \rightarrow \mathbb{R}_{1}^{n+2}$ are non-constant differentiable functions.

Performing calculations similar to those made in Sections 3 and 4, the following equations can be obtained:

$$
\begin{align*}
\lambda_{1} \lambda_{2} a^{\top} & =-\frac{c_{k}}{2}\binom{n}{k+1} \nabla H_{k+1}^{2}-2 c_{k}\left(S \circ P_{k}\right)\left(\nabla H_{k+1}\right)+2 c_{k} P_{k}\left(\nabla H_{k}\right),  \tag{31}\\
\lambda_{1} \lambda_{2}\langle N, a\rangle & =c_{k} L_{k}\left(H_{k+1}\right)-\left(\operatorname{tr}\left(S^{2} \circ P_{k}\right)-c_{k} H_{k}+\lambda_{1}+\lambda_{2}\right) c_{k} H_{k+1},  \tag{32}\\
\lambda_{1} \lambda_{2}\langle\psi, a\rangle & =c_{k}^{2} H_{k+1}^{2}-\left(c_{k} H_{k}-\lambda_{1}\right)\left(c_{k} H_{k}-\lambda_{2}\right)-c_{k} L_{k}\left(H_{k}\right) \tag{33}
\end{align*}
$$

where $c_{k}=(n-k)\binom{n}{k}=(k+1)\binom{n}{k+1}$.
The following example exhibits hypersurfaces of $L_{k}$-2-type in the hyperbolic space $\mathbb{H}^{n+1}$.

Example 4. For each positive number $r$ and each integer $m, 1 \leq m \leq n-1$, let $M_{m}^{n}(r)$ be the $n$-dimensional submanifold of $\mathbb{R}_{1}^{n+2}$ defined by

$$
M_{m}^{n}(r)=\left\{\left(x_{1}, \ldots, x_{n+2}\right) \mid-x_{1}^{2}+\sum_{i=2}^{m+1} x_{i}^{2}=-1-r^{2}, \sum_{j=m+2}^{n+2} x_{j}^{2}=r^{2}\right\}
$$

It is well known that $M_{m}^{n}(r)$ is a complete and non-compact hypersurface of the hyperbolic space $\mathbb{H}^{n+1}$; in fact, $M_{m}^{n}(r)$ is isometric to the standard Riemannian product $\mathbb{H}^{m}\left(-\sqrt{1+r^{2}}\right) \times \mathbb{S}^{n-m}(r)$.

The Gauss map of $M_{m}^{n}(r)$ in $\mathbb{H}^{n+1}$ is given by

$$
N(x)=\left(\frac{r}{\sqrt{1+r^{2}}} x_{1}, \ldots, \frac{r}{\sqrt{1+r^{2}}} x_{m+1}, \frac{\sqrt{1+r^{2}}}{r} x_{m+2}, \ldots, \frac{\sqrt{1+r^{2}}}{r} x_{n+2}\right)
$$

and then $M_{m}^{n}(r)$ has two constant distinct principal curvatures given by

$$
\kappa_{1}=\cdots=\kappa_{m}=\frac{-r}{\sqrt{1+r^{2}}} \quad \text { and } \quad \kappa_{m+1}=\cdots=\kappa_{n}=\frac{-\sqrt{1+r^{2}}}{r} .
$$

Hence, the mean curvature $H_{k}$ is constant for every $k$.
If we put $\psi_{1}=\left(x_{1}, \ldots, x_{m+1}, 0, \ldots, 0\right)$ and $\psi_{2}=\left(0, \ldots, 0, x_{m+2}, \ldots, x_{n+2}\right)$, then $\psi=\psi_{1}+\psi_{2}$ and, by using (7), we obtain $L_{k} \psi_{1}=\lambda_{1} \psi_{1}$ and $L_{k} \psi_{2}=\lambda_{2} \psi_{2}$, where
$\lambda_{1}=\frac{c_{k}}{\sqrt{1+r^{2}}}\left(r H_{k+1}+\sqrt{1+r^{2}} H_{k}\right) \quad$ and $\quad \lambda_{2}=\frac{c_{k}}{r}\left(\sqrt{1+r^{2}} H_{k+1}+r H_{k}\right)$.
Therefore, $M_{m}^{n}(r)$ is a hypersurface of $L_{k}$-2-type of the hyperbolic space $\mathbb{H}^{n+1}$.
Now, we are ready to prove the following classification result.
Theorem 10. Let $\psi: M^{n} \rightarrow \mathbb{H}^{n+1} \subset \mathbb{R}_{1}^{n+2}$ be an orientable $H_{k}$-hypersurface and assume that $M^{n}$ has at most two distinct principal curvatures. Then $M^{n}$ is of $L_{k}$-2-type if and only if $M^{n}$ is an open portion of $M_{m}^{n}(r)$, for some positive integer $m, 1 \leq m \leq n-1$, and for some positive number $r$.

Proof. Let us assume that $M^{n}$ is a hypersurface of $L_{k}$-2-type. Let $\kappa_{1}$ and $\kappa_{2}$ denote the principal curvatures of $M^{n}$, with multiplicities $m$ and $n-m$, respectively. Consider $\left\{E_{1}, E_{2}, \ldots, E_{n}\right\}$ a local orthonormal frame of principal directions of $S$ such that $S E_{i}=\kappa_{1} E_{i}$, for $i=1, \ldots, m$, and $S E_{j}=\kappa_{2} E_{j}, j=m+1, \ldots, n$. Without loss of generality, we can distinguish two cases according to the multiplicity $m$.

Case 1. $m=1$.
Let us consider the open set

$$
\mathcal{U}_{k+1}=\left\{p \in M^{n} \mid \nabla H_{k+1}^{2}(p) \neq 0\right\} .
$$

Our goal is to show that $\mathcal{U}_{k+1}$ is empty. Otherwise, since $M^{n}$ is a $L_{k}-2$-type hypersurface and the mean curvature $H_{k}$ is constant, by taking covariant derivative in (33) we obtain $\lambda_{1} \lambda_{2} a^{\top}=c_{k}^{2} \nabla H_{k+1}^{2}$, that jointly with (31) yields

$$
\begin{equation*}
\left(S \circ P_{k}\right)\left(\nabla H_{k+1}^{2}\right)=-\frac{c_{k}(2 k+3)}{2(k+1)} H_{k+1} \nabla H_{k+1}^{2} \quad \text { on } \mathcal{U}_{k+1} . \tag{34}
\end{equation*}
$$

From the inductive definition of $P_{k+1}=\binom{n}{k+1} H_{k+1} I-S \circ P_{k}$ and (34) we obtain

$$
\begin{equation*}
P_{k+1}\left(\nabla H_{k+1}^{2}\right)=D_{k} H_{k+1} \nabla H_{k+1}^{2} \quad \text { on } \mathcal{U}_{k+1}, \tag{35}
\end{equation*}
$$

where $D_{k}=\frac{2 k+5}{2}\binom{n}{k+1}$. The vector field $\nabla H_{k+1}^{2}$ can be written as $\nabla H_{k+1}^{2}=$ $\sum_{i=1}^{n}\left\langle\nabla H_{k+1}^{2}, E_{i}\right\rangle E_{i}$, and then we get

$$
P_{k+1}\left(\nabla H_{k+1}^{2}\right)=\sum_{i=1}^{n}\left\langle\nabla H_{k+1}, E_{i}\right\rangle \mu_{k+1}^{i} E_{i} .
$$

Hence, Eq. (35) is equivalent to

$$
\left\langle\nabla H_{k+1}^{2}, E_{i}\right\rangle\left(\mu_{k+1}^{i}-D_{k} H_{k+1}\right)=0 \quad \text { on } \mathcal{U}_{k+1},
$$

for every $i=1, \ldots, n$. Therefore, for every $i$ such that $\left\langle\nabla H_{k+1}^{2}, E_{i}\right\rangle \neq 0$ we get

$$
\mu_{k+1}^{i}=D_{k} H_{k+1} .
$$

We will distinguish two cases: (a) $\left\langle\nabla H_{k+1}^{2}, E_{1}\right\rangle \neq 0$, and (b) $\left\langle\nabla H_{k+1}^{2}, E_{i}\right\rangle \neq 0$ for some $i>1$.
(a) First, let us suppose that $\left\langle\nabla H_{k+1}^{2}, E_{1}\right\rangle \neq 0$. Then, we get

$$
\mu_{k+1}^{1}=D_{k} H_{k+1}=\frac{2 k+5}{2} \mu_{k+1}=\frac{2 k+5}{2}\left(\kappa_{1} \mu_{k}^{1}+\mu_{k+1}^{1}\right) .
$$

This equation, bearing in mind that $\binom{n}{k} H_{k}=\mu_{k}=\kappa_{1} \mu_{k-1}^{1}+\mu_{k}^{1}$, leads to

$$
\begin{equation*}
\left.-(2 k+3) \mu_{k+1}^{1} \mu_{k-1}^{1}=(2 k+5)\binom{n}{k} H_{k}-\mu_{k}^{1}\right) \mu_{k}^{1} . \tag{36}
\end{equation*}
$$

Now, by using that $\mu_{j}^{1}=\binom{n-1}{j} \kappa_{2}^{j}$ for $j \in\{1, \ldots, n-1\}$, we can rewrite (36) as follows

$$
A \kappa_{2}^{k}+B=0
$$

where $A$ and $B$ are two nonzero constants. Therefore, $\kappa_{2}$ is constant. This implies, since $H_{k}$ is constant, that the principal curvature $\kappa_{1}$ is constant, and so $H_{k+1}$ is also constant, which is a contradiction.
(b) Now, suppose that $\left\langle\nabla H_{k+1}^{2}, E_{i}\right\rangle \neq 0$ for some $i>1$. Then, we get

$$
\kappa_{1} \mu_{k}^{1, i}+\mu_{k+1}^{1, i}=\mu_{k+1}^{i}=D_{k} H_{k+1}=\frac{2 k+5}{2}\left(\kappa_{1} \mu_{k}^{1}+\mu_{k+1}^{1}\right) .
$$

It is not difficult to see that this equation is equivalent to

$$
\binom{n-2}{k} \kappa_{1}+\binom{n-2}{k+1} \kappa_{2}=\frac{2 k+5}{2}\left(\binom{n-1}{k} \kappa_{1}+\binom{n-1}{k+1} \kappa_{2}\right) .
$$

In other words, $C \kappa_{1}=D \kappa_{2}$, where $C$ and $D$ are two nonzero constants given by

$$
C=\frac{3-n(2 k+3)}{2(n-1)}\binom{n-1}{k} \quad \text { and } \quad D=\frac{n(2 k+3)-1}{2(n-1)}\binom{n-1}{k+1} .
$$

By direct computation, we find that

$$
C\binom{n}{k} H_{k}=\left[\binom{n-1}{k} C+\binom{n-1}{k-1} D\right] \kappa_{2}^{k} .
$$

Therefore, $\kappa_{2}$ is constant. As before, this implies that the $(k+1)$-th mean curvature $H_{k+1}$ is also constant, which is not possible.

Case 2. $1<m<n-1$ (i.e. the multiplicities of two principal curvatures are greater than one).

Without loss of generality, suppose that $\kappa_{1}, \kappa_{2} \neq 0$. By using a standard reasoning involving the Codazzi equations, we deduce that $E_{i}\left(\kappa_{1}\right)=0$, for $i=1, \ldots, m$, and $E_{j}\left(\kappa_{2}\right)=0$, for $j=m+1, \ldots, n$. Since the number of distinct principal curvatures is two, the distribution corresponding to each principal curvature is smooth and integrable (see, e.g., [2, Paragraph 16.10] and [23]). Hence, we deduce that each principal curvature $\kappa_{i}$ is constant on each integral submanifold of the corresponding distribution of the space of principal vectors $V\left(\kappa_{i}\right)$ (see [23]). Therefore, $M^{n}$ is locally isometric to the Riemannian product $M_{1} \times M_{2}$, where $M_{i}$ is the maximal integral submanifold corresponding to the distribution of the space $V\left(\kappa_{i}\right)$ (see, e.g., [18, p. 182]).

Since $H_{k}$ is constant on the hypersurface $M_{1} \times M_{2}$ and $\kappa_{1}$ is constant on $M_{1}$, we deduce that $\kappa_{2}$ is also constant on $M_{1}$. Similarly, the constancy of $H_{k}$ and $\kappa_{2}$ on $M_{2}$ implies that $\kappa_{1}$ is also constant on $M_{2}$. In other words, the principal curvatures $\kappa_{1}$ and $\kappa_{2}$ are constant on the whole hypersurface, and so $H_{k+1}$ is also constant, which is a contradiction.

In conclusion, the mean curvatures $H_{k}$ and $H_{k+1}$ of the hypersurface $M^{n}$ are constant. Since $M^{n}$ has at most two distinct principal curvatures, we get that $M^{n}$ is an isoparametric hypersurface of the hyperbolic space. Bearing in mind the classification of isoparametric hypersurfaces in $\mathbb{H}^{n+1}$ (see [3]), we deduce that $M^{n}$ is an open portion of $M_{m}^{n}(r)$, for some positive integer $m, 1 \leq m \leq n-1$, and for some positive number $r$.

In the case $k=n-1$ we can drop the condition on the principal curvatures of the hypersurface $M^{n}$.

Theorem 11. Let $\psi: M^{n} \rightarrow \mathbb{H}^{n+1} \subset \mathbb{R}_{1}^{n+2}$ be an orientable $H_{n-1}$-hypersurface. If $M^{n}$ is of $L_{n-1}$-2-type then its Gauss-Kronecker curvature $H_{n}$ is a nonzero constant.

Proof. Let us suppose that $H_{n}$ is non constant and consider the nonempty open set

$$
\mathcal{U}_{n}=\left\{p \in M^{n} \mid \nabla H_{n}^{2}(p) \neq 0\right\} .
$$

Since $M^{n}$ is of $L_{n-1}-2$-type and $H_{n-1}$ is constant, by taking covariant derivative in (33) we have $\lambda_{1} \lambda_{2} a^{\top}=c_{n-1}^{2} \nabla H_{n}^{2}$, and by putting this into Eq. (31) we obtain

$$
\begin{equation*}
\left(S \circ P_{n-1}\right)\left(\nabla H_{n}^{2}\right)=-\frac{2 n+1}{2} H_{n} \nabla H_{n}^{2} \quad \text { on } \mathcal{U}_{n} . \tag{37}
\end{equation*}
$$

Since $P_{n}=0$, we deduce $S \circ P_{n-1}=H_{n} I$, and so $S \circ P_{n-1}\left(\nabla H_{n}^{2}\right)=H_{n} \nabla H_{n}^{2}$, that jointly with (37) implies $H_{n} \nabla H_{n}^{2}=0$ on $\mathcal{U}_{n}$, which can not be possible. Therefore, $H_{n}$ is constant and nonzero.

## References

1. L. J. Alías and N. Gürbüz, An extension of Takahashi theorem for the linearized operators of the higher order mean curvatures, Geom. Dedicata, 121 (2006), 113-127.
2. A. L. Besse, Einstein Manifolds, Ergeb. Math. Grenzgeb. (3), 10. Springer-Verlag, Berlin, 1987.
3. E. Cartan, Familles de surfaces isoparamétriques dans les espaces a courbure constante, Ann. Mat. Pura Appl., 17 (1938), 177-191.
4. E. Cartan, Sur des familles remarquables d'hypersurfaces isoparamétriques dans les espaces sphériques, Math. Z., 45 (1939), 335-367.
5. E. Cartan, Sur Quelque Familles Remarquables d'hypersurfaces, C. R. Congrès Math. Liège, 1939, 30-41.
6. E. Cartan, Sur des familles d'hypersurfaces isoparamétriques des espaces sphériques à 5 et à 9 dimensions, Univ. Nac. Tucumán Rev. Ser. A, 1 (1940), 5-22.
7. B. Y. Chen, Total Mean Curvature and Submanifolds of Finite Type, Series in Pure Mathematics, 1. World Scientific Publishing Co., Singapore, 1984.
8. B. Y. Chen, Finite Type Submanifolds and Generalizations, University of Rome, Rome, 1985.
9. B. Y. Chen, Finite type submanifolds in pseudo-Euclidean spaces and applications, Kodai Math. J., 8 (1985), 358-375.
10. B. Y. Chen, Finite-type pseudo-Riemannian submanifolds, Tamkang J. Math., 17(2) (1986), 137-151.
11. B. Y. Chen, 2-type submanifolds and their applications, Chinese J. Math., 14 (1986), 1-14.
12. B. Y. Chen, M. Barros and O. J. Garay, Spherical finite type hypersurfaces, Algebras Groups Geom., 4 (1987), 58-72.
13. B. Y. Chen, Submanifolds of finite type in hyperbolic spaces, Chinese J. Math., 20(1) (1992), 5-21.
14. B. Y. Chen, A report on submanifolds of finite type, Soochow J. Math., 22(2) (1996), 117-337.
15. B. Y. Chen, Some open problems and conjectures on submanifolds of finite type: recent development, Tamkang J. Math., 45(1) (2014), 87-108.
16. T. Hasanis and T. Vlachos, Spherical 2-type hypersurfaces, J. Geom., 40 (1991), 82-94.
17. S. M. B. Kashani, On some $L_{1}$-finite type (hyper)surfaces in $\mathbb{R}^{n+1}$, Bull. Korean Math. Soc., 46 (2009), 35-43.
18. S. Kobayashi and K. Nomizu, Foundations of Differential Geometry, Vol. I, WileyInterscience, New York, NY, USA, 1963; Vol. II, 1969.
19. P. Lucas and H. F. Ramírez-Ospina, Hypersurfaces in pseudo-Euclidean spaces satisfying a linear condition on the linearized operator of a higher order mean curvature, Differential Geom. Appl., 31 (2013), 175-189.
20. P. Lucas and H. F. Ramírez-Ospina, $L_{k}$-2-type Hypersurfaces in $\mathbb{S}^{4}$, submitted for publication.
21. A. Mohammadpouri and S. M. B. Kashani, On some $L_{k}$-finite-type Euclidean hypersurfaces, ISRN Geom., Vol. 2012, article ID 591296, 23 pages.
22. B. O'Neill, Semi-Riemannian Geometry with Applications to Relativity, Academic Press, 1983, New York, London.
23. T. Otsuki, Minimal hypersurfaces in a Riemannian manifold of constant curvature, Amer. J. Math., 92(1) (1970), 145-173.
24. R. Reilly, Variational properties of functions of the mean curvatures for hypersurfaces in space forms, J. Differential Geom., 8 (1973), 465-477.
25. B. Segre, Famiglie di ipersuperficie isoparametrische negli spazi euclidei ad un qualunque numero di demesioni, Atti Accad. Naz. Lincei Rend. Cl. Sci. Fis. Mat. Natur., 27 (1938), 203-207.
26. C. Somigliana, Sulle relazione fra il principio di Huygens e l'ottica geometrica, Atti Accad. Sci. Torino Cl. Sci. Fis. Mat. Natur., 54 (1918-1919), 974-979.
D. Pascual Lucas Saorín

Departamento de Matemáticas
Universidad de Murcia
Campus de Espinardo
30100 Murcia SPAIN
E-mail: plucas@um.es
D. Hector-Fabián Ramírez-Ospina

Departamento de Matemáticas
Universidad de Murcia
Campus de Espinardo
30100 Murcia
Spain
E-mail: hectorfabian.ramirez@um.es


[^0]:    Received March 29, 2014, accepted May 21, 2014.
    Communicated by Bang-Yen Chen.
    2010 Mathematics Subject Classification: 53C40, 53B25.
    Key words and phrases: Hyperbolic hypersurface, Linearized operator $L_{k} ; L_{k}$-finite-type hypersurface, Higher order mean curvatures, Newton transformations.
    This work has been partially supported by MINECO (Ministerio de Economía y Competitividad) and FEDER (Fondo Europeo de Desarrollo Regional) Project MTM2012-34037.
    Second author is supported by FPI Grant BES-2010-036829.
    *Corresponding author.

