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## DeSitter distances in Hilbert spaces

*Dedicated to Zenon Moszner  
on the occasion of his 70th birthday, in friendship*

**Abstract.** All 2-point invariants of an arbitrary DeSitter manifold are determined without any regularity assumption, especially those which are additive on geodesics.

### 1. Introduction

Two functional equations play a role in this note: the functional equation of *2-point invariants* of a DeSitter manifold will be solved, and, moreover, the functional equation of *additivity along lines* of such a manifold will be studied within the set of solutions of the first equation. This then leads us to the notion of *distance* in a DeSitter World over an arbitrary pre-Hilbert space  $X$  of dimension at least 3. In the case  $\dim X < \infty$  the problems above are treated in chapter 4 of the book [1]. However, further methods and ideas are needed in the present situation in comparison with the finite-dimensional case.

For the classical theory of DeSitter's World see [3], [4], for modern developments see [5], [6], [7].

### 2. Points, motions, lines

Let  $X$  be a real pre-Hilbert space, i.e. a real vector space furnished with an inner product

$$\delta : X \times X \rightarrow \mathbb{R}, \quad \delta(x, y) =: xy,$$

satisfying  $x^2 > 0$  for all  $x \neq 0$  in  $X$ . We assume that the dimension of  $X$  is at least 3. Let  $t$  be a fixed element of  $X$  such that  $t^2 = 1$ . Define

$$H := t^\perp := \{x \in X \mid tx = 0\}.$$

Then  $H \oplus \mathbb{R}t = X$  holds true (see section 2 of [2]). Since  $x - (tx)t$  is in  $H$  for all  $x \in X$ ,

$$x = \bar{x} + x_0t \quad \text{with } \bar{x} \in H \quad \text{and } x_0 \in \mathbb{R}$$

implies  $\bar{x} = x - (tx)t$  and  $x_0 = tx$ . Define

$$S(X) := \{x \in X \mid \bar{x}^2 - x_0^2 = 1\}$$

as the set of *points* of *DeSitter's Manifold* over  $X$ . The Lorentz-Minkowski distance  $l(x, y)$  of  $x, y \in X$  is defined (see [2]) by the expression

$$l(x, y) = (\bar{x} - \bar{y})^2 - (x_0 - y_0)^2.$$

A Lorentz transformation of  $X$  is a mapping  $\lambda$  from  $X$  into itself such that

$$l(x, y) = l(\lambda(x), \lambda(y))$$

holds true for all  $x, y \in X$ . In [2] all Lorentz transformations of  $X$  are determined.

The restriction of a surjective Lorentz transformation  $\lambda : X \rightarrow X$  on  $S(X)$  will be called a *motion* of  $S(X)$  provided that  $\lambda(0) = 0$ . Motions of  $S(X)$  must be bijective. We hence get the group  $\Delta(X)$  of motions of  $S(X)$ .

The points  $a, b$  of  $S(X)$  are called *separated* if, and only if,  $b \notin \{a, -a\}$ . Such a pair must be linearly independent. Otherwise an equation  $\alpha a = \beta b$  holds true with real  $\alpha, \beta$  which are not both 0. But this implies

$$\alpha^2 = l(\alpha a, 0) = l(\beta b, 0) = \beta^2,$$

i.e.  $\beta = \alpha$  or  $\beta = -\alpha$ .

Suppose that  $a, b$  are points of  $S(X)$  which are separated. Then every ellipse, every euclidean line, every branch of a hyperbola in

$$\{\xi a + \eta b \mid \xi, \eta \in \mathbb{R}\} \cap S(X) \tag{1}$$

is called a *line* of  $S(X)$ . All  $\xi a + \eta b$  in (1) are characterized by the equation

$$(\xi + \gamma\eta)^2 + (1 - \gamma^2)\eta^2 = 1, \tag{2}$$

where  $\gamma := \gamma(a, b) := \bar{a}\bar{b} - a_0b_0$  designates the *pseudo-euclidean scalar product* of  $a, b$ . In the cases  $\gamma^2 < 1$ ,  $\gamma^2 = 1$ ,  $\gamma^2 > 1$ , respectively, we get an ellipse (a *closed line*), two euclidean lines (*null-lines*), two branches of a hyperbola (*open lines*), respectively, of  $S(X)$ . The lines of  $S(X)$  are also called its *geodesics*.

### 3. The functional equation of 2-point invariants

Suppose that  $W \neq \emptyset$  is a set and that

$$d : S(X) \times S(X) \rightarrow W$$

satisfies

$$d(x, y) = d(f(x), f(y)) \quad (3)$$

for all motions  $f$  of  $S(X)$  and all  $x, y \in S(X)$ . Then  $d$  is called a *2-point invariant* of  $S(X)$ .

#### THEOREM 1

$\Delta(X)$  acts transitively on  $S(X)$ . If  $a, b$  and  $c, e$  are pairs of separated points, then there exists  $\delta \in \Delta(X)$  with  $\delta(a) = c$  and  $\delta(b) = e$  if, and only if,

$$\gamma(a, b) = \gamma(c, e)$$

holds true.

*Proof.* a) In step c) of the proof of Theorem 3 in [2] we showed that to  $x, y \in X \setminus \{0\}$  there exists a bijective Lorentz transformation  $\lambda$  with  $\lambda(0) = 0$  and  $\lambda(x) = y$  if, and only if,  $l(x, 0) = l(y, 0)$ . Suppose that  $x, y$  are points of  $S(X)$ . Then

$$l(x, 0) = 1 = l(y, 0)$$

holds true. There hence exists a motion  $\delta$  with  $\delta(x) = y$ .

b) If  $a, b$  and  $c, e$  are pairs of separated points, and if  $\delta \in \Delta(X)$  satisfies  $\delta(a) = c$  and  $\delta(b) = e$ , then

$$l(a, b) = l(c, e),$$

since  $\delta$  is a Lorentz transformation. Hence

$$(\bar{a} - \bar{b})^2 - (a_0 - b_0)^2 = (\bar{c} - \bar{e})^2 - (c_0 - e_0)^2,$$

i.e.  $\gamma(a, b) = \gamma(c, e)$ , in view of  $a, b, c, e \in S(X)$ .

c) Let  $a, b$  and  $c, e$  be pairs of separated points satisfying  $\gamma(a, b) = \gamma(c, e)$ . Because of Theorem 1 in [2], Lorentz transformations of  $X$  fixing 0, must be linear. Separated points  $x, y$  must thus be transformed into separated points under motions. In view of step a) we hence may assume  $a = c$  without loss of generality. by observing

$$\gamma(x, y) = \gamma(f(x), f(y)) \quad (4)$$

for  $x, y \in S(X)$  and motions  $f$ . The situation now is that  $a, b$  and  $a, e$  are pairs of separated points such that  $\gamma(a, b) = \gamma(a, e)$ . If  $h \in H$  satisfies  $h^2 = 1$ , which especially implies  $h \in C(X)$ , then we even may assume that  $a = h$ , in

view of step a). Then  $\gamma(a, b) = \gamma(a, e)$  reads  $h\bar{b} = h\bar{e}$ , i.e.  $hb = he$  because of  $h \perp t$ . We now consider the pre-Hilbert space

$$X_0 := \{x - (xh)h \mid x \in X\}.$$

Obviously,  $t \in X_0$ . Again, we would like to apply Theorem 3 of [2], but this time for  $X_0$  and for the points

$$\xi := b - (bh)h \quad \text{and} \quad \eta := e - (eh)h.$$

Observe

$$\begin{aligned} l(\xi, 0) &= \bar{\xi}^2 - b_0^2 = (\bar{b} - (bh)h)^2 - b_0^2 \\ &= 1 - (bh)^2 = 1 - (eh)^2 \\ &= l(\eta, 0). \end{aligned}$$

There hence exists a Lorentz transformation  $\lambda_0$  of  $X_0$  satisfying  $\lambda_0(0) = 0$  and

$$\lambda_0(b - (bh)h) = e - (eh)h. \quad (5)$$

The problem now is to extend  $\lambda_0$  to a Lorentz transformation  $\lambda$  of  $X$  by putting

$$\lambda(x) := \lambda_0(x - (xh)h) + (xh)h$$

for all  $x \in X$ . That  $\lambda$  is an extension of  $\lambda_0$  follows from

$$xh = 0 \quad \text{for all } x \in X_0.$$

Put  $x_h := x - (xh)h$  for  $x \in X$ . Then

$$\overline{\lambda(x)} = \overline{\lambda_0(x_h)} + (xh)h =: x_1 + (xh)h.$$

Put  $[\lambda_0(x_h)]_0 =: x_2$ . Then

$$\begin{aligned} l(\lambda(x), \lambda(y)) &= (x_1 + (xh)h - y_1 - (yh)h)^2 - (x_2 - y_2)^2 \\ &= l(\lambda_0(x_h), \lambda_0(y_h)) + A \end{aligned}$$

with  $A = 2(x_1 - y_1)h(xh - yh) + (xh - yh)^2 = (xh - yh)^2$  since

$$x_1 = \overline{\lambda_0(x_h)} \in X_0$$

implies  $x_1h = 0$ . Similarly, we get

$$l(x, y) = l(x_h, y_h) + A.$$

i.e.  $l(x, y) = l(\lambda_0(x_h), \lambda_0(y_h)) + A = l(\lambda(x), \lambda(y))$ , and  $\lambda$  must hence be a Lorentz transformation of  $X$ . We finally would like to show

$$\lambda(h) = h \quad \text{and} \quad \lambda(b) = e.$$

In fact,

$$\lambda(h) = \lambda_0(h - h^2 \cdot h) + h^2 h = \lambda_0(0) + h = h,$$

and

$$\lambda(b) = \lambda_0(b - (bh)h) + (bh)h = [e - (eh)h] + (eh)h = e,$$

in view of (5) and  $bh = eh$ .

We now would like to solve the functional equation (3) of 2-point invariants.

#### THEOREM 2

Let  $g : \mathbb{R} \rightarrow W$  be a function and let  $w_0, w_1$  be fixed elements of  $W$ . Then

$$d(x, y) = \begin{cases} g(\bar{x}\bar{y} - x_0y_0) & \text{for } x, y \text{ separated} \\ w_0 & \text{for } x = y \\ w_1 & \text{for } x = -y \end{cases} \quad (6)$$

with  $x, y \in S(X)$  is a solution of (3). If, on the other hand,

$$d : S(X) \times S(X) \rightarrow W$$

solves (3), then there exists a function  $g : \mathbb{R} \rightarrow W$  and elements  $w_0, w_1 \in W$  such that (6) holds true.

*Proof.* Obviously, (6) solves (3) for all motions  $f$  and all  $x, y \in S(X)$ , in view of (4). Assume now that  $d : S(X) \times S(X) \rightarrow W$  is a solution of (3). Take elements  $i, j \in H$  with  $i^2 = 1, j^2 = 1, ij = 0$ . For  $k \in \mathbb{R}$  define  $g(k)$  by means of

$$g(k) := d(i, ki + j + kt). \quad (7)$$

Observe here  $i \in S(X), ki + j + kt \in S(X)$  and

$$\gamma(i, ki + j + kt) = k. \quad (8)$$

Moreover, put  $w_0 := d(i, i)$  and  $w_1 := d(i, -i)$ . If  $x \in S(X)$ , there exists  $f \in \Delta(X)$  with  $f(i) = x$  on account of Theorem 1. Hence

$$d(x, x) = d(f(i), f(i)) = d(\bar{i}, \bar{i}) = w_0.$$

Since  $f$  is linear, we also get

$$d(x, -x) = d(f(i), f(-i)) = d(i, -i) = w_1.$$

Suppose now that  $x, y \in S(X)$  are separated. If  $\gamma(x, y) =: k$ , then, according to (8) and Theorem 1, there exists  $f \in \Delta(X)$  satisfying

$$f(i) = x \quad \text{and} \quad f(ki + j + kt) = y.$$

Hence  $d(x, y) = d(f(i), f(ki + j + kt)) = d(i, ki + j + kt) = g(k)$ , in view of (7). Thus

$$d(x, y) = g(k) = g(\gamma(x, y)) = g(\bar{x}\bar{y} - x_0y_0).$$

#### 4. The additivity equation

Since motions are linear, images of null-lines must be null-lines, and images of closed lines, open lines must be closed lines, open lines, respectively.

Suppose that  $s : S(X) \times S(X) \rightarrow \mathbb{R}_{\geq 0} := \{r \in \mathbb{R} \mid r \geq 0\}$  satisfies the following property.

(A) *If  $x, y$  are separated points on a line  $l$ , and if  $z \in l$  is between  $x, y$ , then*

$$s(x, y) = s(x, z) + s(z, y) \quad (9)$$

*holds true.*

Then  $s$  is called *additive*.

If  $l$  is null or open, the usual betweenness relation is meant. If  $l$  is closed, then  $z$  is assumed to be an element of the smaller part of the ellipse in question of the two parts defined by the points  $x, y \in l$ .

#### THEOREM 3

*If  $s : S(X) \times S(X) \rightarrow \mathbb{R}_{\geq 0}$  is an additive 2-point invariant, then there exist non-negative constants  $r_1, r_2$  such that*

$$s(x, y) = 0, \quad s(x, y) = r_1 \operatorname{ch}^{-1}(\gamma(x, y)), \quad s(x, y) = r_2 \cos^{-1}(\gamma(x, y)),$$

*provided  $x, y$  are separated points on a null-line, an open line, a closed line, respectively.*

Here  $\operatorname{ch}^{-1} \alpha = \beta$  with  $\alpha \geq 1$  and  $\beta \geq 0$  is defined by  $\operatorname{ch} \beta = \alpha$ . Similarly,  $\cos^{-1} \alpha = \beta$  with  $\beta \in [0, \pi]$  and  $\alpha \in [-1, 1]$  means  $\cos \beta = \alpha$ .

*Proof.* a) Let  $a + \mathbb{R}v := \{a + \lambda v \mid \lambda \in \mathbb{R}\}$ ,  $v \neq 0$ , be a null-line. Hence

$$\gamma(a, a) = 1, \quad \gamma(a, v) = 0 = \gamma(v, v). \quad (10)$$

Any two distinct points of this line  $l$  must be separated. Moreover,

$$\gamma(a + \lambda_1 v, a + \lambda_2 v) = \gamma(a + \lambda_3 v, a + \lambda_4 v) \quad (11)$$

in the case  $\lambda_1 \neq \lambda_2$  and  $\lambda_3 \neq \lambda_4$ , in view of (10).

Suppose that  $x \neq y$  are on  $l$  and that  $z \notin \{x, y\}$  is on  $l$  between  $x$  and  $y$ . Then (11) and Theorem 1 imply the existence of motions  $\delta_1, \delta_2$  satisfying

$$\delta_1(x) = x, \quad \delta_1(y) = z, \quad \delta_2(x) = z, \quad \delta_2(y) = y.$$

Hence  $s(x, y) = s(x, z)$  and  $s(x, y) = s(z, y)$ . Thus  $s(x, y) = 0$ , in view of (9).

b) Let (1) be an open line satisfying (2) with  $\gamma(a, b) =: k$ . Up to a motion we may assume  $a = i$  and  $b = ki + j + kt$  with  $i, j \in H$ ,  $i^2 = 1$ ,  $j^2 = 1$ ,  $ij = 0$ . The two branches of (2) are given by

$$\operatorname{ch} \lambda \cdot i + \operatorname{sh} \lambda \cdot v \quad \text{and} \quad \operatorname{ch} \lambda \cdot (-i) + \operatorname{sh} \lambda \cdot v$$

with  $\lambda \in \mathbb{R}$  and

$$v := \frac{j + kt}{\sqrt{k^2 - 1}}.$$

Without loss of generality we may work with the first version, since in the other case  $-i$  can be replaced by  $i$ . We hence get

$$p(\lambda) := \operatorname{ch} \lambda \cdot i + \operatorname{sh} \lambda \cdot v, \quad \lambda \in \mathbb{R}, \quad (12)$$

with  $\gamma(i, i) = 1$ ,  $\gamma(i, v) = 0$ ,  $\gamma(v, v) = -1$ . Two points  $p(\lambda_1), p(\lambda_2)$  with  $\lambda_1 \neq \lambda_2$  are separated. Moreover, in this case

$$\gamma(p(\lambda_1), p(\lambda_2)) = \operatorname{ch}(\lambda_1 - \lambda_2) > 1 \quad (13)$$

holds true. ( $k > 1$  need not to be true: if  $a, b$  in (1) are on different branches of the underlying hyperbola, then  $k < -1$ .)

Let  $x, y$  be different points on (12). Because of  $\gamma(x, y) = \gamma(y, x)$  and Theorem 1 there exists a motion that interchanges  $x$  and  $y$ . Hence  $s(x, y) = s(y, x)$ . We hence may assume  $x = p(\lambda_1)$  and  $y = p(\lambda_2)$  with  $\lambda_1 < \lambda_2$ . Theorem 2 implies

$$s(x, y) = g(\operatorname{ch}(\lambda_1 - \lambda_2)) \geq 0,$$

in view of (13). Put  $\varphi(\lambda) = g(\operatorname{ch} \lambda)$  for  $\lambda \geq 0$ . On account of (9) we get for  $\lambda_2 \leq \lambda_3 \leq \lambda_1$ .

$$\varphi(\lambda_1 - \lambda_2) = \varphi(\lambda_1 - \lambda_3) + \varphi(\lambda_3 - \lambda_2).$$

A standard procedure now leads to  $\varphi(\lambda) = r_1 \lambda$  with a non-negative constant  $r_1$ . Hence

$$s(x, y) = g(\operatorname{ch}(\lambda_1 - \lambda_2)) = r_1(\lambda_1 - \lambda_2) = r_1 \operatorname{ch}^{-1}(\gamma(x, y)),$$

on account of (13).

c) Let (1) be a closed line satisfying (2) with  $\gamma(a, b) =: k$ . Mutatis mutandis to the case b), we may work with the line  $l$ ,

$$p(\lambda) = \cos \lambda \cdot i + \sin \lambda \cdot v, \quad \lambda \in \mathbb{R}, \quad (14)$$

with

$$v := \frac{j + kt}{\sqrt{1 - k^2}}.$$

Hence  $\gamma(i, i) = 1$ ,  $\gamma(i, v) = 0$ ,  $\gamma(v, v) = 1$ . Take separated points  $x, y$  on the line  $l$  with

$$x = p(\lambda_1) \quad \text{and} \quad y = p(\lambda_2),$$

$\lambda_2 < \lambda_1$ ,  $\lambda_1 - \lambda_2 < \pi$ . Then

$$\gamma(p(\lambda_1), p(\lambda_2)) = \cos(\lambda_1 - \lambda_2) \in ]-1, 1[.$$

Theorem 2 implies  $s(x, y) = g(\cos(\lambda_1 - \lambda_2)) \geq 0$ . Put  $\varphi(\lambda) = g(\cos \lambda)$  for  $\lambda \in [0, \pi]$ . Again (9) together with a standard procedure leads to  $\varphi(\lambda) = r_2 \lambda$  with a non-negative constant  $r_2$ . Hence

$$s(x, y) = g(\cos(\lambda_1 - \lambda_2)) = r_2(\lambda_1 - \lambda_2) = r_2 \cos^{-1}(\gamma(x, y)).$$

In the case  $r_1 = r_2 = 1$  we call  $s(x, y)$  *DeSitter's distance function*.

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