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DISTRIBUTIONS OF STABLE RANDOM FIELDS OF CHENTSOV TYPE

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§1. Introduction

In this paper we discuss the determinism of distributions of some stable random fields which are constructed through integral-geometric method. The determinism depends on the dimension of the parameter space R^{d} .

We say that a family of random variables $\{X(t); t \in \mathbb{R}^d\}$ is a symmetric α -stable (abbreviated to $S \alpha S$) random field on \mathbb{R}^d if every finite linear combination $\sum_{i=1}^n a_i X(t_i)$ has a symmetric stable distribution of index α . Let (E, \mathscr{B}, μ) be a measure space. We say that a family of random variables $\{Y(B); B \in \mathscr{B}, \mu(B) < \infty\}$ is the $S \alpha S$ random measure corresponding to (E, \mathscr{B}, μ) if (i) $E(\exp[izY(B)]) = \exp(-\mu(B)|z|^{\alpha})$, for $z \in \mathbb{R}$ and $\mu(B) < \infty$, (ii) $Y(B_1), Y(B_2), \cdots$ are independent whenever B_1, B_2, \cdots are disjoint and $\mu(B_j) < \infty, j = 1, 2, \cdots$, (iii) $Y(\bigcup_{j=1}^{\infty} B_j) = \sum_{j=1}^{\infty} Y(B_j)$ a.s. whenever B_1, B_2, \cdots are disjoint and $\mu(\bigcup_{j=1}^{\infty} B_j) < \infty$.

We define a class of $S \alpha S$ random fields with a particular choice of E. Let E_0 be the set of all (d - 1)-dimensional spheres in \mathbf{R}^d . Any element of E_0 is expressed by a coordinate system (r, x), where (r, x) corresponds to the sphere with radius $r \in \mathbf{R}_+ = (0, \infty)$ and center $x \in \mathbf{R}^d$. Thus we make the identification

(1.1)
$$E_0 = \{(r, x); r \in \mathbf{R}_+, x \in \mathbf{R}^i\}.$$

For $t \in \mathbf{R}^{d}$, let s_{t} be the set of all spheres which separate the point t and the origin O, namely

(1.2)
$$S_t = \{(r, x); d(x, O) \leq r\} \triangle \{(r, x); d(r, x) \leq r\},\$$

where $A \triangle B$ denotes the symmetric difference of A and B and d(a, b) denotes the Euclidean distance between a and b. Let \mathscr{B}_0 be the σ -algebra

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of Borel sets in E_0 . Given a measure μ on (E_0, \mathscr{B}_0) such that

(1.3)
$$\mu(S_t) < \infty$$
 for all $t \in \mathbf{R}^d$,

we define an $S\alpha S$ random field by

(1.4)
$$X(t) = Y(S_t), \quad t \in \mathbf{R}^d$$

using the $S \alpha S$ random measure $\{Y(B)\}$ corresponding to $(E_0, \mathscr{B}_0, \mu)$. We call this $\{X(t)\}$ S αS random field of Chentsov type on \mathbf{R}^d associated with μ .

Such a random field is viewed as an extension of N.N. Chentsov's representation $Y(S'_t)$ of Lévy's Brownian motion of \mathbf{R}^{d} -parameter. The $Y(S'_t)$ is defined by Chentsov through Gaussian random measure Y associated with a measure on the space E' of all hyperplanes of co-dimension 1 in \mathbf{R}^{d} and the defining set S'_t is the set of all hyperplanes which separate t and the origin O, [1], [3]. S. Takenaka, [7], applied this idea to stable case. Using E_0 in place of E', he proves that if $d\mu_{\beta}(r, x) = r^{\beta - a - 1} dr dx$, $0 < \beta < 1$, then the Chentsov type $S \alpha S$ random field $X_{\alpha,\beta}(t)$ associated with (E_0, μ_{β}) is self-similar with exponent $H = \beta/\alpha$. For d = 1, $\{X_{\alpha,\beta}(t)\}$ presents a new example of $S \alpha S$, H-self-similar process with stationary increments in the area of α and H where no examples were known before.

The distributions of a Chentsov type $S\alpha S$ random field on \mathbb{R}^d have a characteristic property which depends on the dimension d of the parameters. We do not assume any condition other than (1.3) for the associated measure. The aim of this paper is to prove the following theorem.

THEOREM 1. Let $0 < \alpha < 2$. Let μ be a measure on (E_0, \mathscr{B}_0) satisfying (1.3) and let $\{X(t); t \in \mathbf{R}^a\}$ be the S α S random field of Chentsov type on \mathbf{R}^a associated with μ . Then, for any n > d + 1 and for any distinct $t_1, \dots, t_n \in \mathbf{R}^a$, the distribution $(X(t_1), \dots, X(t_n))$ is determined by its (d + 1)-dimensional marginal distributions.

COROLLARY. Let $0 < \alpha < 2$. Let μ and $\tilde{\mu}$ be measures on (E_0, \mathscr{B}_0) satisfying (1.3). Let $\{X(t); t \in \mathbf{R}^a\}$ and $\{\tilde{X}(t); t \in \mathbf{R}^a\}$ be the S α S random fields of Chentsov type associated with μ and $\tilde{\mu}$, respectively. If the (d + 1)dimensional distributions of $\{X(t)\}$ and $\{\tilde{X}(t)\}$ coincide, then $\{X(t)\}$ and $\{\tilde{X}(t)\}$ are equivalent, that is, the finite-dimensional distributions of $\{X(t)\}$ and $\{\tilde{X}(t)\}$ coincide.

Remark 1. The number d + 1 in Theorem 1 is best possible in the following sense. There are two Chentsov type random fields $\{X(t)\}$ and

 $\{\tilde{X}(t)\}$ associated with μ and $\tilde{\mu}$, respectively, such that, for some $T = (t_1, \dots, t_{d+1})$, the d-dimensional marginal distributions of $(X(t_1), \dots, X(t_{d+1}))$ and $(\tilde{X}(t_1), \dots, \tilde{X}(t_{d+1}))$ coincide but their (d + 1)-dimensional distributions are different. (see Example 4.2)

Remark 2. If we take E' and S'_t instead of E_0 and S_t and define

$$X'(t) = Y(S'_t) \quad \text{for } t \in \mathbf{R}^d$$
,

where Y is an $S\alpha S$ random measure with $0 < \alpha < 2$ associated with a measure μ' on E' satisfying $\mu'(S'_t), < \infty$, then we have determinism by d-dimensional marginal distributions instead of determinism by (d + 1)-dimensional marginal distributions in Theorem 1. Namely, for any n > d and any distinct $t_1, \dots, t_n \in \mathbb{R}^d$, the distribution of $(X'(t_1), \dots, X'(t_n))$ is determined by its d-dimensional marginal distributions. This fact can be proved by a similar method as Theorem 1.

Theorem 1 will be reduced to a geometric theorem concerning an intersection property of a family of cones in $\mathbf{R}_+ \times \mathbf{R}^d$. The proof of this geometric theorem is an essential part of our argument. For $t \in \mathbf{R}^d$, set

(1.5)
$$C_t = \{(r, x); d(x, t) \leq r\}$$

Then, C_t is a right cone in $\mathbf{R}_+ \times \mathbf{R}^d$ with vertex (0, t). Note that the point (0, t) is not included in the space $\mathbf{R}_+ \times \mathbf{R}^d$. Hereafter we simply call C_t the cone with vertex t. The relation

$$S_t = C_0 \triangle C_t$$

shows that, instead of S_i 's, we may study C_i 's. Given *m* cones $C_{\iota_1}, \dots, C_{\iota_m}$, we consider the partition of the set $\bigcup_{i=1}^m C_{\iota_i}$ generated by $\{C_{\iota_i}, i = 1, \dots, m\}$. Now set

(1.6)
$$\mathscr{E}_m = \{ e = (e_1, \dots, e_m); e_i = 0 \text{ or } 1 \text{ for } i = 1, \dots, m \} \setminus \{ (0, \dots, 0) \}.$$

We call $e \in \mathscr{E}_m$ a label of size m and \mathscr{E}_m the label set. With the notation

$$C^{\scriptscriptstyle 1}_t = C_t \quad ext{and} \quad C^{\scriptscriptstyle 0}_t = C^c_t = (\mathbf{R}_+ imes \mathbf{R}^d) ackslash C_t \ ,$$

we define

(1.7)
$$C(T, e) = \bigcap_{i=1}^{m} C_{i_i}^{e_i}$$

for $T = (t_1, \dots, t_m) \in (\mathbf{R}^d)^m$ and $e = (e_1, \dots, e_m) \in \mathscr{E}_m$. Then C(T, e), $e \in \mathscr{E}_m$, are disjoint sets and

(1.8)
$$\bigcup_{i=1}^{m} C_{t_i} = \bigcup_{e \in \mathcal{S}_m} C(T, e)$$

For $e = (e_1, \dots, e_m) \in \mathscr{E}_m$, the complementary label e^* of e is defined by

$$e^* = (e_1^*, \dots, e_m^*), \quad e_i + e_i^* = 1 \text{ for any } i.$$

THEOREM 2. If $m \ge d + 3$, then, for any $T = (t_1, \dots, t_m) \in (\mathbf{R}^d)^m$, there exists a label $e \in \mathscr{E}_m$ such that $C(T, e) = \emptyset$ and $C(T, e^*) = \emptyset$.

In §2 we reduce Theorem 1 to Theorem 2, which is proved in §3 after a series of lemmas. Concluding remarks are given in §4. The particular cases d = 1 and 2 of Theorem 1 have been treated in the author's paper [4], [5]. Determinism under different defining sets in one-dimensional case will be discussed in a joint paper with S. Takenaka [6].

§ 2. Reduction of Theorem 1 to Theorem 2

Let $0 < \alpha < 2$. Let μ be a measure on (E_0, \mathscr{B}_0) satisfying (1.3) and let $\{X(t); t \in \mathbf{R}^d\}$ be the $S \alpha S$ random field of Chentsov type on \mathbf{R}^d associated with μ . For $t \in \mathbf{R}^d$ let S_t be the set defined by (1.2). For $T = (t_1, \dots, t_n)$ $\in (\mathbf{R}^d)^n$ and $e = (e_1, \dots, e_n) \in \mathscr{E}_n$, we write

(2.1)
$$S_{l_k}^{e_k} = \begin{cases} S_{l_k} & \text{ if } e_k = 1 \\ S_{l_k}^c & \text{ if } e_k = 0 \end{cases}$$

(2.2)
$$\widetilde{S}_{l_k}^{e_k} = \begin{cases} S_{l_k} & \text{if } e_k = 1 \\ \mathbf{R}_+ \times \mathbf{R}^d & \text{if } e_k = 0 \end{cases}$$

$$(2.3) S(T, e) = \bigcap_{k=1}^{n} S_{i_k}^{e_k}$$

(2.4)
$$\tilde{S}(T, e) = \bigcap_{k=1}^{n} \tilde{S}_{\iota_{k}}^{e_{k}}.$$

S(T, e) is an element (labelled with e) of the partition of the set $\bigcup_{k=1}^{n} S_{\iota_k}$ generated by $\{S_{\iota_k}; k = 1, \dots, n\}$.

DEFINITION 2.1. We say that $T = (t_1, \dots, t_n) \in (\mathbf{R}^d)^n$ satisfies Condition (I) if there exists a label $e \in \mathscr{E}_n$ such that $S(T, e) = \emptyset$.

Remark. Suppose that $T = (t_1, \dots, t_n)$ satisfies Condition (I). Then $T' = (t_1, \dots, t_n, t_{n+1}, \dots, t_m)$ satisfies Condition (I) for any $t_{n+1}, \dots, t_m \in \mathbf{R}^d$. In fact, suppose that $S(T, e) = \emptyset$ for $e = (e_1, \dots, e_n) \in \mathscr{E}_n$. If we take a label $\tilde{e} = (\tilde{e}_1, \dots, \tilde{e}_m)$ such that $\tilde{e}_1 = e_1, \dots, \tilde{e}_n = e_n$, then $S(T', \tilde{e}) = \emptyset$.

LEMMA 2.2. Let t_1, \dots, t_n be different points in \mathbb{R}^d . If $T = (t_1, \dots, t_n)$

satisfies Condition (I), then the distribution of $X_T = (X(t_1), \dots, X(t_n))$ is determined by the system of (n-1)-dimensional marginal distributions of X_T .

Proof. The characteristic function of X_T is written, for $z = (z_1, \dots, z_n) \in \mathbf{R}^n$, as

(2.5)

$$\varphi_{T}(z) = E \exp\left\{i\sum_{k=1}^{n} z_{k}Y(S_{t_{k}})\right\}$$

$$= E \exp\left\{i\sum_{\substack{k=1\\e\in \mathcal{I}_{n}}} z_{k}\sum_{\substack{e\in \mathcal{I}_{n}\\e_{k}=1}} Y(S(T, e))\right\}$$

$$= E \exp\left\{i\sum_{\substack{e\in \mathcal{I}_{n}}} \left(\sum_{k=1}^{n} e_{k}z_{k}\right)Y(S(T, e))\right\}$$

$$= \exp\left\{-\sum_{\substack{e\in \mathcal{I}_{n}}} \left|\sum_{k=1}^{n} e_{k}z_{k}\right|^{\alpha}\mu(S(T, e))\right\}.$$

$$= \exp\left\{-\sum_{\substack{e\in \mathcal{I}_{n}}} \left|\xi(e)\cdot z\right|^{\alpha}(\sharp e)^{\alpha/2}\mu(S(T, e))\right\}.$$

where $e = (e_1, \dots, e_n), \ #e = \sum_{i=1}^n e_i$, and

(2.6)
$$\xi(e) = (1/(\#e)^{1/2})e$$

On the other hand it is known that the characteristic function of an *n*-dimensional $S\alpha S$ distribution, $0 < \alpha < 2$, has the following unique canonical representation:

(2.7)
$$\varphi(z) = \exp\left\{-c\int_{S^{n-1}} |\xi \cdot z|^{\alpha} \lambda(d\xi)\right\},$$

where c > 0 and λ is a symmetric probability measure on the unit sphere S^{n-1} [2]. Comparing the last expression of (2.5) to (2.7) and noticing that $\xi(e)$ of (2.6) belongs to S^{n-1} for any e, we see that the last expression of (2.5) gives the canonical form of $\varphi_T(z)$. So, we see that $\varphi_T(z)$ is determined by the values of $\mu(S(T, e)), e \in \mathscr{E}_n$, and that, conversely, $\mu(S(T, e)), e \in \mathscr{E}_n$, are determined by $\varphi_T(z)$.

Since μ is a measure, the following consistency condition holds:

(2.8)
$$\mu(\tilde{S}(T, e)) = \sum_{e' \in \mathscr{E}_n'(e)} \mu(S(T, e')) \quad \text{for each } e \in \mathscr{E}_n,$$

where, for $e = (e_1, \dots, e_n) \in \mathscr{E}_n$, $\mathscr{E}'_n(e)$ is the subset of \mathscr{E}_n defined by

(2.9)
$$\mathscr{E}'_n(e) = \{ e' = (e'_1, \dots, e'_n); e'_k \ge e_k \text{ for } k = 1, \dots, n \}.$$

Since the number of the elements of \mathscr{E}_n is $2^n - 1$, (2.8) consists of $2^n - 1$ equations. But one of them, that is, the case $e = (1, \dots, 1)$, is a trivial

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equation. So we consider

(2.10)
$$\sum_{e' \in \mathscr{E}_n \setminus \{e\}} \mu(S(T, e')) = \mu(\tilde{S}(T, e)) \quad \text{for } e \in \mathscr{E}_n \setminus \{(1, \cdots, 1)\}$$

as a system of $2^n - 2$ simultaneous linear equations in which the unknowns are the values of $\mu(S(T, e))$, $e \in \mathscr{E}_n$, (the number of them is $2^n - 1$) and the given right-hand sides are the values of $\mu(\tilde{S}(T, e))$, $e \in \mathscr{E}_n \setminus \{(1, \dots, 1)\}$. These right-hand sides are determined by the (n - 1)-dimensional marginal distributions of X_T . We write the matrix representation of the system (2.10) as

$$(2.11) M_n \mathbf{x} = \mathbb{D},$$

where M_n is a $(2^n - 2) \times (2^n - 1)$ -matrix, whose components are 0 or 1. Let $M_n(k)$ be the $(2^n - 2) \times (2^n - 2)$ -matrix obtained from M_n by deleting the k-th column. Then we can prove that

(2.12) $M_n(k)$ is invertible for any $k = 1, \dots, 2^n - 1$.

Proof of (2.12) will be given at the end of this section. By the assumption that T satisfies Condition (I), there exists a label e such that $S(T, e) = \emptyset$. This implies $\mu(S(T, e)) = 0$ for the label e. So, the number of the unknowns is reduced to $2^n - 2$. Suppose that the $\mu(S(T, e))$ corresponding to the label e is the k-th component of the column vector of the unknowns. Then our system of simultaneous linear equations is equivalent to the system having $M_n(k)$ as its coefficient matrix. By virtue of (2.12) the system of equations has a unique solution. Thus, all the unknowns are determined.

Now we need to study the problem when T satisfies Condition (I).

Let $T = (t_1, \dots, t_n) \in (\mathbf{R}^d)^n$ and $e = (e_1, \dots, e_n) \in \mathscr{E}_n$. The set S(T, e) is partitioned into two disjoint subsets:

$$(2.13) S(T, e) = \{S(T, e) \cap C_0\} \cup \{S(T, e) \cap C_0^e\},$$

where C_0 is the cone with vertex 0.

LEMMA 2.3. We have

(2.14)
$$S(T, e) \cap C_0 = C(T, e^*) \cap C_0,$$

(2.15) $S(T, e) \cap C_0^c = C(T, e) \cap C_0^c,$

where e^{*} is the complementary element of e. Proof. We have

$$S(T,e) \cap \ C_{\scriptscriptstyle 0} = \left(igcap_{i=1}^n S^{e_i}_{t_i}
ight) \cap \ C_{\scriptscriptstyle 0} = igcap_{i=1}^n \left(S^{e_i}_{t_i} \cap \ C_{\scriptscriptstyle 0}
ight).$$

If $e_i = 1$, then

$$S^{e_i}_{t_i} \cap \, C_{\scriptscriptstyle 0} = S_{t_i} \cap \, C_{\scriptscriptstyle 0} = (C_{t_i} riangle \, C_{\scriptscriptstyle 0}) \cap \, C_{\scriptscriptstyle 0} = C^{\, c}_{t_i} \cap \, C_{\scriptscriptstyle 0} = C^{e_i *}_{t_i} \cap \, C_{\scriptscriptstyle 0} \, .$$

If $e_i = 0$, then

$$S_{\iota_i}^{e_i} \cap C_0 = S_{\iota_i}^c \cap C_0 = (C_{\iota_i} \triangle C_0)^c \cap C_0 = C_{\iota_i} \cap C_0 = C_{\iota_i}^{e_i^*} \cap C_0.$$

Hence

$$\bigcap_{i=1}^{n} (S_{t_{i}}^{e_{i}} \cap C_{0}) = \bigcap_{i=1}^{n} (C_{t_{i}}^{e_{i}*} \cap C_{0}) = \left(\bigcap_{i=1}^{n} C_{t_{i}}^{e_{i}*}\right) \cap C_{0} = C(T, e^{*}) \cap C_{0}.$$

This proves (2.14). The relation (2.15) is obtained more easily.

Using Lemma 2.3 we can reduce Condition (I) to the following Condition (II).

DEFINITION 2.4. We say that $T = (t_1, t_2, \dots, t_m) \in (\mathbf{R}^d)^m$ satisfies Condition (II) if there exists a label $e \in \mathscr{E}_m$ such that both $C(T, e) = \emptyset$ and $C(T, e^*) = \emptyset$ hold.

LEMMA 2.5. $T = (t_1, \dots, t_n) \in (\mathbf{R}^d)^n$ satisfies Condition (I) if and only if $\tilde{T} = (0, t_1, \dots, t_n) \in (\mathbf{R}^d)^{n+1}$ satisfies Condition (II).

Proof. (i) Suppose that $T = (t_1, \dots, t_n)$ satisfies Condition (I). Let $e = (e_1, \dots, e_n) \in \mathscr{E}_n$ be the label such that $S(T, e) = \emptyset$. It follows from (2.13) and Lemma 2.3 that $C(T, e^*) \cap C_0 = \emptyset$ and $C(T, e) \cap C_0^c = \emptyset$. Put $\tilde{e} = (0, e_1, \dots, e_n) \in \mathscr{E}_{n+1}$. Then

$$egin{aligned} C(T, e^*) \, \cap \, C_{\scriptscriptstyle 0} &= \, C(T, ilde{e}^*) \,, \ C(T, e) \, \cap \, C_{\scriptscriptstyle 0}^c &= \, C(ilde{T}, ilde{e}) \,. \end{aligned}$$

Hence \tilde{T} satisfies Condition (II).

(ii) Suppose that $\tilde{T} = (0, t_1, \dots, t_n)$ satisfies Condition (II). Let $\tilde{e} = (e_0, e_1, \dots, e_n) \in \mathscr{E}_{n+1}$ be the label such that $C(\tilde{T}, \tilde{e}) = \emptyset$ and $C(\tilde{T}, \tilde{e}^*) = \emptyset$. Let $e = (e_1, \dots, e_n)$. If $e_0 = 0$, then $e_0^* = 1$ and

$$egin{aligned} & arnothing & = C(T, ilde{e}) \, = \, C(T,e) \, \cap \, C^e_{_0} \, , \ & arnothing & = \, C(ilde{T}, ilde{e}^*) \, = \, C(T,e^*) \, \cap \, C_e \, , \end{aligned}$$

which implies $S(T, e) = \emptyset$ by (2.13) and Lemma 2.3. If $e_0 = 1$, then $e_0^* = 0$ and, taking account of $(e^*)^* = e$, we have $S(T, e^*) = \emptyset$. In either case, $T = (t_1, \dots, t_n)$ satisfies Condition (I).

Proof that Theorem 2 implies Theorem 1. We assume that Theorem 2 is valid. Let $n \ge d + 1$. Let t_1, \dots, t_n be distinct points in \mathbb{R}^d . Theorem 2 combined with Lemma 2.5 tells us that $T = (t_1, \dots, t_n)$ satisfies Condition (I). Hence, by Lemma 2.2 the distribution of X_T is determined by its (n-1)-dimensional marginal distributions. Further, if $n-1 \ge d+2$, then the (n-1)-dimensional marginal distributions of X_T are determined by their (n-2)-dimensional marginal distributions. Proceeding in this way we see that the distributions of X_T is determined by its (d+1)dimensional marginal distributions. Theorem 1 is proved. \Box

Proof of (2.12). To write down the matrix M_n , we introduce a linear order among the element of \mathscr{E}_n . Let $n(e) = \sum_{i=1}^n 2^{i-1}e_i$ for $e = (e_1, \dots, e_n)$. We define e < e' if n(e) < n(e'). This gives a linear order in \mathscr{E}_n . Thus the first element is $(1, 0, 0, \dots, 0)$ and the last element is $(1, \dots, 1)$. We have

$$M_{2} = \begin{pmatrix} 1 & 0 & 1 \\ 0 & 1 & 1 \end{pmatrix}$$
$$M_{3} = \begin{pmatrix} 1 & 0 & 1 & 0 & 1 & 0 & 1 \\ 0 & 1 & 1 & 0 & 0 & 1 & 1 \\ 0 & 0 & 1 & 0 & 0 & 1 & 1 \\ 0 & 0 & 0 & 1 & 1 & 1 & 1 \\ 0 & 0 & 0 & 0 & 1 & 0 & 1 \\ 0 & 0 & 0 & 0 & 0 & 1 & 1 \end{pmatrix}$$

Let $N_n = 2^n - 1 = 2N_{n-1} + 1$. Then M_n is a $(N_n - 1) \times N_n$ -matrix. It is easy to see that

for $n \ge 3$. Let a_j be the *j*-th column vector and a_{ij} be the (i, j)-component of M_n . If we delete the last column in M_n , then we get an upper triangular matrix $M_n(N_n)$ with diagonal elements 1. Hence

are linearly independent. Now we claim the following.

(2.18) If
$$c_1\mathfrak{a}_1 + c_2\mathfrak{a}_2 + \cdots + c_{N_n}\mathfrak{a}_{N_n} = 0$$
 with $c_k = 0$
for some $k \neq N_n$, then $c_{N_n} = 0$.

Suppose that (2.18) is true. Then we see that $M_n(k)$ is invertible for every k. Indeed, if $k = N_n$, then $M_n(k)$ is invertible by (2.17). If $k \neq N_n$, then (2.17) and (2.18) show that the column vectors of $M_n(k)$ are linearly independent.

It remains to prove the assertion (2.18). It suffices to show that the relation

(2.19)
$$c_1 a_1 + c_2 a_2 + \cdots + c_{N_n} a_{N_n} = 0$$
 with $c_{N_n} = 1$

implies that

(2.20)
$$c_k \neq 0 \text{ for } k = 1, \dots, N_n - 1.$$

Note that, by (2.17), all of c_1, \dots, c_{N_n-1} are determined uniquely by (2.19). Denote the row vector

$$\mathbb{C}_n = (c_1, c_2, \cdots, c_{N_n})$$

Using the column vectors of M_{n-1} in place of those of M_n , we get the row vector \mathbb{C}_{n-1} in place of \mathbb{C}_n . Let

 $\mathbb{C}_{n-1} = (\gamma_1, \gamma_2, \cdots, \gamma_{N_{n-1}}), \qquad \gamma_{N_{n-1}} = 1.$

We write (2.19) componentwise:

(2.21)
$$c_1a_{i1} + c_2a_{i2} + \cdots + c_{N_n}a_{iN_n} = 0, i = 1, \dots, N_n - 1, \text{ and } c_{N_n} = 1.$$

For $i = N_{n-1} + 2, \dots, N_n - 1$, the relation between M_n and M_{n-1} in (2.16) shows that (2.21) implies

(2.22)
$$(c_{N_{n-1}+2}, \cdots, c_{N_n}) = (\gamma_1, \cdots, \gamma_{N_{n-1}}).$$

For $i = N_{n-1}$, (2.21) reduces to $c_{N_{n-1}} + c_{N_n} = 0$ by virtue of (2.16). Hence

$$(2.23) c_{N_{n-1}} = -1.$$

For $i = 1, \dots, N_{n-1} - 1$, taking account of (2.16) and using (2.22) and (2.23), we have

(2.24)
$$(c_1, c_2, \cdots, c_{N_{n-1}}) = (-\gamma_1, -\gamma_2, \cdots, -\gamma_{N_{n-1}})$$

It follows from (2.22) and (2.24) that

(2.25) $C_n = (-C_{n-1}, C_{N_{n-1}+1}, C_{n-1}).$

We have

$$(2.26) C_2 = (-1, -1, 1)$$

explicitly from M_2 . For $i = N_{n-1} + 1$, (2.21) reduces to

 $c_{N_{n-1}+1} + c_{N_{n-1}+2} + \cdots + c_{N_n} = 0$

by (2.16). Hence, noticing (2.22) and using (2.26) or (2.25) for n-1 in place of n, we get

(2.27)
$$c_{N_{n-1}+1} = \begin{cases} 1 & \text{for } n = 3 \\ -\gamma_{N_{n-2}+1} & \text{for } n \ge 4 \end{cases}.$$

Now, from (2.25) and (2.27) we see that each component of c_n is 1 or -1. This proves (2.20).

§3. Proof of Theorem 2

We prepare lemmas.

LEMMA 3.1. Let $t_i \in \mathbf{R}^d$, $i = 1, \dots, n + m$. Let $A = \{t_1, \dots, t_n\}$ and $B = \{t_{n+1}, \dots, t_{n+m}\}$.

$$(3.1) \qquad \qquad \bigcap_{t_i \in A} C_{t_i} \subset \bigcup_{t_j \in B} C_{t_j}$$

if and only if

(3.2)
$$\max_{t_i \in A} d(t_i, x) \ge \min_{t_j \in B} d(t_j, x) \quad \text{for any } x \in \mathbf{R}^d.$$

We denote the relation (3.1) by A < B. This means

$$(\bigcap_{t_i\in A}C_{t_i})\cap (\bigcap_{t_j\in B}C_{t_j}^c)=\varnothing$$
.

Proof. Suppose (3.1). Let $x \in \mathbb{R}^d$. Let $r = \max_{t_i \in A} d(x, t_i)$. Then $d(x, t_i) \leq r$ for any $t_i \in A$, that is, $(r, x) \in \bigcap_{t_i \in A} C_{t_i}$. Hence $(r, x) \in \bigcup_{t_j \in B} C_{t_j}$. This means $d(x, t_j) \leq r$ for some $t_j \in B$. Hence (3.2) holds. Conversely, assume (3.2). Let $(r, x) \in \bigcap_{t_i \in A} C_{t_i}$. That means $d(x, t_i) \leq r$ for every $t_i \in A$.

It follows from (3.2) that there exists t_{j_0} such that $d(t_{j_0}, x) \leq r$. So, $(r, x) \in C_{t_{j_0}}$.

LEMMA 3.2. Let $1 \leq k \leq d+1$. Let t_1, \dots, t_k , and t_{k+1}, \dots, t_{d+2} in \mathbb{R}^d be such that there is no hyperplane of co-dimension 1 containing t_1, \dots, t_k , t_{k+1}, \dots, t_{d+1} . Suppose that there exist positive constants p_1, \dots, p_k and q_{k+1}, \dots, q_{d+2} such that

(3.3)
$$\sum_{i=1}^{k} p_i = \sum_{j=k+1}^{d+2} q_j = 1,$$

(3.4)
$$\sum_{i=1}^{k} p_i t_i = \sum_{j=k+1}^{d+2} q_j t_j.$$

Then at least one of the following holds:

(3.5)
$$\bigcap_{i=1}^{k} C_{i_i} \subset \bigcup_{j=k+1}^{d+2} C_{i_j} \quad (\text{that is } A \prec B),$$

$$(3.6) \qquad \qquad \bigcup_{i=1}^k C_{t_i} \supset \bigcap_{j=k+1}^{d+2} C_{t_j} \quad (that \ is \ A \succ B) \,,$$

where $A = \{t_1, \dots, t_k\}$ and $B = \{t_{k+1}, \dots, t_{d+2}\}$.

Proof. Let D be the (d-1)-dimensional sphere on which the points $t_1, \dots, t_k, t_{k+1}, \dots, t_{d+1}$ lie. Without loss of generality, we assume that the center of D is $O = (0, \dots, 0) \in \mathbf{R}^d$. Let r be the radius of D. Suppose that $|t_{d+2}| \leq r$. We will show that (3.5) holds. By Lemma 3.1 it is enough to show that, for any $x \in \mathbf{R}^d$,

(3.7)
$$\max_{i=1,\dots,k} d(t_i, x) - \min_{j=k+1,\dots,d+2} d(t_j, x) \ge 0.$$

Taking account of (3.3) and (3.4) we have

(3.8)
$$\min_{i=1,\dots,k}(t_i, x) \leq \sum_{i=1}^k p_i(t_i, x) = \sum_{j=k+1}^{d+2} q_j(t_j, x) \leq \max_{j=k+1,\dots,d+2}(t_j, x)$$

where (x, y) denotes inner product of R^i . Let i_0 and j_0 be the elements which attain the minimum and the maximum in (3.8), respectively. Then

$$(3.9) (t_{j_0} - t_{i_0}, x) \ge 0.$$

On the other hand,

$$(3.10) \qquad \begin{aligned} \max_{i=1,\cdots,k} \{d(t_i, x)\}^2 &- \min_{j=k+1,\cdots,d+2} \{d(t_j, x)\}^2 \\ &\geq \{d(t_{i_0}, x)\}^2 - \{d(t_{j_0}, x)\}^2 \\ &= \{|t_{i_0}|^2 + |x|^2 - 2(t_{i_0}, x)\} - \{|t_{j_0}|^2 + |x|^2 - 2(t_{j_0}, x)\} \\ &= |t_{i_0}|^2 - |t_{j_0}|^2 + 2(t_{j_0} - t_{i_0}, x) \,. \end{aligned}$$

(3.9) and the assumption $|t_{i_0}| = r \ge |t_{j_0}|$ imply that the last term of (3.10) is non-negative. So, (3.7) is proved.

Suppose that $|t_{d+2}| \ge r$. Then we prove that

(3.11)
$$\max_{j=k+1,\dots,d+2} d(t_j, x) - \min_{i=1,\dots,k} d(t_i, x) \ge 0 \quad \text{for any } x \in \mathbb{R}^d$$

which implies (3.6) by Lemma 3.1. In fact, let

$$(t_{i_0}, x) = \max_{i=1,\dots,k} (t_i, x), \quad (t_{j_0}, x) = \min_{j=k+1,\dots,d+2} (t_{j_0}, x).$$

Then

$$egin{aligned} &\max_{j=k+1,\cdots,d+2} \{d(t_j,\,x)\}^2 - \min_{i=1,\cdots,k} \{d(t_i,\,x)\}^2 \ &\geq \{d(t_{j_0},\,x)\}^2 - \{d(t_{i_0},\,x)\}^2 \ &= |t_{j_0}|^2 - |t_{i_0}|^2 + 2(t_{i_0} - t_{j_0},\,x) \geq 0 \end{aligned}$$

which is (3.11).

LEMMA 3.3. Let $t_1, \dots, t_{d+1} \in \mathbb{R}^d$. Suppose that no hyperplane of codimension 1 contains them and that any d vectors out of t_1, \dots, t_{d+1} are linearly independent. Let $t_{d+2} = 0$. Then the set $\{t_1, \dots, t_{d+1}, t_{d+2}\}$ is uniquely partitioned into two disjoint sets A, B such that $A \neq \emptyset$, $B \ni t_{d+2}$ and there exist positive constants p_i 's and q_j 's satisfying

(3.12)
$$\sum_{i_i \in A} p_i t_i = \sum_{i_j \in B} q_j t_j,$$

(3.13)
$$\sum_{i_i \in A} p_i = \sum_{i_j \in B} q_j = 1.$$

Proof. Since t_1, \dots, t_{d+1} are linearly dependent, there exist constants c_1, \dots, c_{d+1} such that $(c_1, \dots, c_{d+1}) \neq (0, \dots, 0)$ and $\sum_{i=1}^{d+1} c_i t_i = 0$. Notice that $c_i \neq 0$ for any *i* by the assumption that any *d* out of t_1, \dots, t_{d+1} are linearly independent. Moreover, c_1, \dots, c_{d+1} are unique up to constant multiple. We have $\sum_{i=1}^{d+1} c_i \neq 0$, because, if it is zero, then $\sum_{i=1}^{d} c_i (t_i - t_{d+1}) = 0$ and t_1, \dots, t_{d+1} are on a hyperplane of co-dimension 1. So, we may assume that $\sum_{i=1}^{d+1} c_i > 0$. Let $A = \{t_i; c_i > 0\}$ and $B = \{t_i; c_i < 0\} \cup \{t_{d+2}\}$. Let $p_i = c_i$ for $c_i > 0$, $q_j = -c_j$ for $c_j < 0$, and $q_{d+2} = \sum_{i=1}^{d+1} c_i$. Then $\sum_{t_i \in A} p_i - \sum_{t_j \in B} q_j = 0$ and (3.12) holds. Multiplication of some constant yields (3.13). Uniqueness of A and B is obvious from this argument. \Box

COROLLARY 3.4. Let $t_i \in \mathbb{R}^d$, $i = 1, \dots, d+2$. Assume that no d+1 points out of them are contained in a hyperplane of co-dimension 1 in \mathbb{R}^d .

Then the set $\{t_1, \dots, t_{d+2}\}$ is partitioned into two disjoint non-empty sets A and B such that, for some $p_i > 0$ and $q_j > 0$,

(3.14)
$$\sum_{t_i \in A} p_i t_i = \sum_{t_j \in B} q_j t_j, \quad \sum_{t_i \in A} p_i = \sum_{t_j \in B} q_j = 1.$$

The partition is unique up to the naming of A and B.

Proof. Let $u_i = t_i - t_{d+2}$ and apply Lemma 3.3 to u_1, \dots, u_{d+2} .

We call A, B in Corollary 3.4 the natural partition of $\{t_1, \dots, t_{d+2}\}$.

The corollary above is rephrased geometrically as follows. For a finite set $C = \{t_1, \dots, t_n\} \subset \mathbf{R}^d$, denote by \overline{C} the solid simplex having C as the set of vertices, that is,

$$\overline{C} = \left\{\sum_{i=1}^n p_i t_i; \sum_{i=1}^n p_i = 1, p_i \ge 0, i = 1, \cdots, n\right\}.$$

COROLLARY 3.5. Let t_i , $i = 1, \dots, d + 2$, be as in Corollary 3.4. Then there are two disjoint non-empty sets A, B such that $A \cup B = \{t_1, \dots, t_{d+2}\}$, $A \cap B = \emptyset$, and $\overline{A} \cap \overline{B} \neq \emptyset$. The sets A, B are unique up to naming of A and B. The set $\overline{A} \cap \overline{B}$ consists of only one point.

Combining Lemma 3.1 and Corollary 3.4, we get the following proposition.

PROPOSITION 3.6. For any $T = (t_1, \dots, t_{d+2}) \in (\mathbf{R}^d)^{d+2}$ such that no d + 1 points out of t_1, \dots, t_{d+2} are contained in a hyperplane of co-dimension 1 in \mathbf{R}^d , there exists a label $e \in \mathscr{E}_{d+2}$ which satisfies $C(T, e) = \emptyset$.

Now we deal with a set of d + 3 points in \mathbb{R}^d in order to discuss Condition (II). Consider a set $\Gamma = \{t_1, \dots, t_{d+3}\}$ in \mathbb{R}^d . Assume that Γ is non-degenerate in the sense that

(3.15) no d + 1 points out of t_1, \dots, t_{d+3} are contained in a hyperplane of co-dimension 1 in \mathbf{R}^d .

For each *i*, apply Corollary 3.4 to $\Gamma \setminus \{t_i\}$ and let

$$(3.16) \Gamma \setminus \{t_i\} = A_i \cup B_i$$

be the natural partition of $\Gamma \setminus \{t_i\}$. By Lemma 3.2, at least one of $A_i \prec B_i$ and $A_i \succ B_i$ holds.

Let $i \neq j$. We say that t_i and t_j link together if the restrictions to $\Gamma \setminus \{t_i, t_j\}$ of the natural partitions of $\Gamma \setminus \{t_i\}$ and $\Gamma \setminus \{t_j\}$ coincide.

LEMMA 3.7. Let $i \neq j$ and suppose that t_i and t_j link together. Let A_i , B_i and A_j , B_j be the natural partitions of $\Gamma \setminus \{t_i\}$ and $\Gamma \setminus \{t_j\}$, respectively. If

$$(3.17) A_i \prec B_i, \quad A_j \succ B_j, \quad A_i \cap A_j \neq \emptyset,$$

then $T = (t_1, \dots, t_{d+3})$ satisfies Condition (II).

Proof. Without loss of generality we assume i = 1, j = 2. Keeping $A_1 \cap A_2 \neq \emptyset$ in mind, we can find A and B satisfying $A \cup B = \Gamma \setminus \{t_1, t_2\}$ and $A \cap B = \emptyset$ such that one of the following four conditions holds:

(a) $A_1 = A \cup \{t_2\}, B_1 = B, A_2 = A \cup \{t_1\}, B_2 = B;$ (b) $A_1 = A \cup \{t_2\}, B_1 = B, A_2 = A, B_2 = B \cup \{t_1\};$ (c) $A_1 = A, B_1 = B \cup \{t_2\}, A_2 = A, B_2 = B \cup \{t_1\};$ (d) $A_1 = A, B_1 = B \cup \{t_2\}, A_2 = A \cup \{t_1\}, B_2 = B.$ We may assume that $A = \{3, \dots, k\}$ and $B = \{k + 1, \dots, d + 3\}$ where

 $3 \leq k \leq d+3 \ (B= \oslash \text{ if } k=d+3).$

Case (a). We have

$$C(T, e) = \emptyset \quad \text{with} \quad e = (e_1, \underbrace{1, 1, \cdots, 1}_{k-1}, \underbrace{0, \cdots, 0}_{d+3-k}),$$

$$C(T, e') = \emptyset \quad \text{with} \quad e = (0, e'_2, \underbrace{0, \cdots, 0}_{k-2}, \underbrace{1, \cdots, 1}_{d+3-k}),$$

whatever e_1 and e'_2 are. Letting $e_1 = 1$ and $e'_2 = 0$, we get a complementary pair e, e'. Hence T satisfies Condition (II).

Case (b). We have

$$C(T, e) = \emptyset \quad \text{with} \quad e = (e_1, \underbrace{1, 1, \cdots, 1}_{k-1}, \underbrace{0, \cdots, 0}_{d+3-k})$$
$$C(T, e') = \emptyset \quad \text{with} \quad e' = (1, e'_2, \underbrace{0, \cdots, 0}_{k-2}, \underbrace{1, \cdots, 1}_{d+3-k}),$$

whatever e_1 and e'_2 are. Letting $e_1 = 0$ and $e'_2 = 0$, we obtain a complementary pair.

Cases (c) and (d) are treated similarly to (a) and (b), respectively. \Box

Remark. Another sufficient condition for T to satisfy Condition (II) is that there exists i such that $A_i < B_i$ and $A_i > B_i$. But we will not use this condition.

We see easily that, to prove Theorem 2, it is enough to prove it for m = d + 3. In order to prove it for m = d + 3 under the condition that

 $\{t_1, \dots, t_{d+3}\}$ are non-degenerate in the sense of (3.15), we will show the existence of *i* and *j* which satisfy the condition of Lemma 3.7. Applying Corollary 3.4 to $\Gamma \setminus \{t_1\}$ and $\Gamma \setminus \{t_2\}$, we have

(3.18)
$$\sum_{k=1}^{d+3} c_{1k} t_k = 0$$
 with $c_{11} = 0$, $\sum_{k=1}^{d+3} c_{1k} = 0$, $c_{1k} \neq 0$ $(k \neq 1)$,

and

$$(3.19) \quad \sum_{k=1}^{d+3} c_{2k} t_k = 0 \quad \text{with} \quad c_{22} = 0, \quad \sum_{k=1}^{d+3} c_{2k} = 0, \quad c_{2k} \neq 0 \ (k \neq 2) \ .$$

The representations are unique up to constant multiplication. We assume $c_{12} > 0$ and $c_{21} > 0$. We set, for $i \ge 3$,

(3.20)
$$\begin{cases} \lambda_1 = c_{2i}/c_{1i} \\ c_{ik} = c_{2k} - \lambda_i c_{1k} \end{cases}$$

Then we get the relations for $\Gamma \setminus \{i\}, i = 3, \dots, d + 3$, that

(3.21)
$$\sum_{k=1}^{d+3} c_{ik} t_k = 0 \quad \text{with} \quad c_{ii} = 0, \quad \sum_{k=1}^{d+3} c_{ik} = 0.$$

Obviously we have, for $i \geq 3$,

(3.22)
$$\begin{cases} c_{i1} = c_{21} > 0\\ c_{i2} = -\lambda_i c_{12}\\ c_{ik} = c_{1k} (\lambda_k - \lambda_i) = c_{2k} (1 - \lambda_i / \lambda_k), \quad \text{for } k \ge 3. \end{cases}$$

Moreover we see that λ_i , $i = 3, \dots, d+3$, are distinct and $c_{ik} \neq 0$ for $k \neq i$, because, if otherwise, some d+1 points in Γ are contained in a hyperplane of co-dimension 1. Without loss of generality we assume $\lambda_i < \lambda_{i+1}$ for $i = 3, \dots, d+3$. Let

$$(3.23) L_{-} = \{i \ge 3; \lambda_i < 0\}, L_{+} = \{j \ge 3; \lambda_j > 0\}.$$

We see that $c_{i2} > 0$ for $i \in L_{-}$ and $c_{j2} < 0$ for $j \in L_{+}$. Using the relations in (3.22) and noticing that the natural partition of $\Gamma \setminus \{i\}$ is made according to the signs of c_{ik} , we get

LEMMA 3.8. If both i and i + 1 belong to L_{-} , then t_i and t_{i+1} link together. If both j and j + 1 belong to L_{+} , then t_j and t_{j+1} link together.

Now we assume that

$$(3.24) L_{-} \neq \emptyset \quad \text{and} \quad L_{+} \neq \emptyset$$

The case without this assumption will be treated later. Let

 $L_{-} = \{3, 4, \cdots, r\}, \qquad L_{+} = \{r + 1, \cdots, d + 3\}.$

Then we get the following lemma.

LEMMA 3.9. The following pairs link together: (1) t_1 and t_3 ; (2) t_1 and t_{d+3} ; (3) t_2 and t_7 ; (4) t_2 and t_{r+1} .

Proof. Again use (3.22) and the fact that the natural partition of $\Gamma \setminus \{t_i\}$ is decided by the signs of c_{ik} , $k \neq i$.

It follows from Lemma 3.7 that, if L_{-} or L_{+} contains adjacent elements i, j satisfying (3.17), then $T = (t_1, \dots, t_{d+3})$ satisfies Condition (II). So, let us consider the situation that neither L_{-} nor L_{+} contains adjacent elements satisfying (3.17). In the naming of A_i , B_i in the natural partition (3.16) of $\Gamma \setminus \{t_i\}$, we make $A_i \ni t_1$ for $i = 2, 3, \dots, d+3$, and $A_1 \ni t_2$. Recalling that the natural partitions are made by the signs of c_{ik} , we see that $t_2 \in A_i$ for $i \in L_{-}$ and that $t_2 \in B_i$ for $i \in L_{+}$. We note that $A_{i-1} \cap A_i \neq \emptyset$ for $i \in L_{-} \cup L_{+}$. Hence Lemma 3.7 yields that we have one of the following situations:

- (1) $A_i \prec B_i$ for $i \in L_- \cup L_+$;
- (2) $A_i \prec B_i$ for $i \in L_-$ and $A_j \succ B_j$ for $j \in L_+$;
- (3) $A_i > B_i$ for $i \in L_- \cup L_+$;
- (4) $A_i > B_i$ for $i \in L_-$ and $A_j \prec B_j$ for $j \in L_+$.

We will prove that in each case at least one of pairs (1), (2), (3), (4) of Lemma 3.9 satisfies the condition of Lemma 3.7.

Case (1). If $A_1 \prec B_1$, then t_1 and t_{d+3} satisfy the condition of Lemma 3.7, because $A_{d+3} \prec B_{d+3}$ and $A_1 \cap B_{d+3} \ni t_2$. If $A_1 \succ B_1$, then t_1 and t_3 satisfy (3.17), since $A_3 \prec B_3$ and $A_1 \cap A_3 \ni t_2$.

Case (2). If $A_2 \prec B_2$, then t_2 and t_{r+1} satisfy (3.17), since $A_{r+1} \succ B_{r+1}$ and $A_2 \cap A_{r+1} \ni t_1$. If $A_2 \succ B_2$, then t_2 and t_r satisfy (3.17), because $A_r \prec B_r$ and $A_2 \cap A_r \ni t_1$.

Case (3). Similarly to case (1), the pair t_1 , t_3 or the pair t_1 , t_{d+3} satisfies the condition of Lemma 3.7.

Case (4). Similar to case (2). The pair t_2 , t_7 or the pair t_2 , t_{7+1} satisfies (3.17).

Thus, under the assumption (3.24), $T = (t_1, \dots, t_{d+3})$ satisfies Condition (II).

Let us consider the case where L_{-} or L_{+} is empty.

LEMMA 3.10. If $L_{-} = \emptyset$, then each of the following pairs links together:

 t_1, t_2 ; t_1, t_{d+3} ; t_2, t_3 .

If $L_{+} = \emptyset$, then each of the following pairs links together.

$$t_1, t_2$$
; t_1, t_3 ; t_2, t_{d+3} .

Proof. Suppose that $L_{-} = \emptyset$. Let

$$A = \{i \ge 3; c_{1i} > 0, c_{2i} > 0\}, \qquad B = \{i \ge 3; c_{1i} < 0, c_{2i} < 0\}.$$

Then $A \cup B = \{3, \dots, d+3\}$, and hence t_1 and t_2 link together. If $L_+ = \emptyset$, then letting

$$A = \{i \geq 3; c_{1i} > 0, c_{2i} < 0\}, \qquad B = \{i \geq 3; c_{1i} < 0, c_{2i} > 0\},$$

we see that $A \cup B = \{3, \dots, d+3\}$ and that t_1 and t_2 link together. The other assertions are proved in the same way by use of (3.22).

As before we make the naming of A_i , B_i in the natural partition (3.16) in such a way that $A_i \ni t_1$ for $i = 2, 3, \dots, d+3$, and $A_1 \ni t_2$. We have $t_2 \in A_i$ for $i \in L_-$ and $t_2 \in B_i$ for $i \in L_+$.

Suppose that $L_{-} = \emptyset$. If L_{+} contains adjacent elements i, j satisfying (3.17), then, by Lemmas 3.7 and 3.8, $T = (t_1, \dots, t_{d+3})$ satisfies Condition (II). So, suppose that L_{+} does not contain adjacent elements satisfying (3.17). Then we have one of the following:

(1) $A_i \prec B_i ext{ for } i \geq 3$, (2) $A_i \succ B_i ext{ for } i \geq 3$.

Case (1). If $A_1 \prec B_1$, then t_1 , t_{d+3} satisfy condition of Lemma 3.7 since $A_{d+3} \prec B_{d+3}$ and $A_1 \cap B_{d+3} \ni t_2$. If $A_2 \succ B_2$, then t_2 , t_3 satisfy (3.15), since $A_3 \prec B_3$ and $A_2 \cap A_3 \ni t_1$. In the remaining case, suppose that $A_1 \succ B_1$ and $A_2 \prec B_2$. If $c_{1k} > 0$ for some $k \ge 3$, then $c_{2k} > 0$ and $A_1 \cap A_2 \ni t_k$. If $c_{1k} < 0$ for some $k \ge 3$, then $c_{2k} < 0$ and $B_1 \cap B_2 \ni t_k$. So, t_1 , t_2 satisfy the condition of Lemma 3.7. We made use of Lemma 1.10.

Case (2). If $A_1 > B_1$, then t_1 , t_{d+3} satisfy the condition of Lemma 3.7,

since $A_{d+3} > B_{d+3}$ and $A_1 \cap B_{d+3} \ni t_2$. If $A_2 \prec B_2$, then t_2 , t_3 satisfy the condition, because $A_3 > B_3$ and $A_2 \cap A_3 \ni t_1$. If $A_1 \prec B_1$ and $A_2 > B_2$, then t_1 , t_2 satisfy the condition by same reason as case (1).

Suppose that $L_{+} = \emptyset$. Then we can make similar discussion. Namely, suppose that L_{-} does not contain adjacent elements satisfying (3.15). Then (1) or (2) holds. In either case we can find the following pair satisfying the condition of Lemma 3.7.

Case (1). If $A_1 > B_1$, then t_1 , t_2 . If $A_2 > B_2$, then t_2 , t_{d+3} . If $A_1 < B_1$ and $A_2 < B_2$, then t_1 , t_2 .

Case (2). If $A_1 \prec B_1$, then t_1, t_3 . If $A_2 \prec B_2$, then t_2, t_{d+3} . If $A_1 \succ B_1$ and $A_2 \succ B_2$, then t_1, t_2 .

Therefore, in the case that L_{-} or L_{+} is empty, $T = (t_1, \dots, t_{d+3})$ satisfies Condition (II). This finishes proof of Theorem 2 for m = d + 3 under the assumption that t_1, \dots, t_{d+3} are non-degenerate in the sense of (3.15).

If d + 1 points are on a hyperplane of co-dimension 1 and no d + 2 points are on a hyperplane of co-dimension 1, then we can apply Lemma 3.2 again and similar argument can be made. If d + 2 points are on a hyperplane of co-dimension 1, then, taking account of the remark to Definition 2.1, we see that the situation is reduced to (d - 1)-dimensional case.

§4. Concluding remarks

In order to construct an example mentioned in Remark 1 of §1, we prepare a lemma.

LEMMA 4.1. Let $T = (t_1, \dots, t_{d+2}) \in (\mathbb{R}^d)^{d+2}$, where t_1, \dots, t_{d+2} are distinct and no d + 1 points of them are on a hyperplane of codimension 1. Let Dbe the (d-1)-dimensional sphere on which the points t_1, \dots, t_{d+1} lie. Assume that t_{d+2} is situated inside of D and, moreover, that $\overline{A} \cap \overline{B} \neq \emptyset$ for $A = \{t_{d+1}, t_{d+2}\}$ and $B = \{t_1, \dots, t_d\}$, using the notation introduced before Corollary 3.5. Then there is no label e of size d + 2 such that $C(T, e) = C(T, e^*) = \emptyset$.

Proof. For $e = (e_1, \dots, e_{d+2}) \in \mathscr{E}_{d+2}$, let $A_e = \{t_i; e_i = 1\}$ and $B_e = \{t_i; e_i = 0\}$. In order to prove our assertion, it is enough to consider only e such that $A_e \ni t_{d+2}$. We separate our discussion into three cases.

(a) A_e and B_e give the natural partition of $\{t_1, \dots, t_{d+2}\}$.

- (b) Either A_e or B_e is a one point set.
- (c) The remaining case.

Case (a). We have $A_e = A$ and $B_e = B$ by the assumption. From the proof of Lemma 3.2 we see that A > B. We do not have A < B. In fact, we can find a (d-1)-dimensional sphere D' such that $D' \supset B$ and that the points t_{d+1} , t_{d+2} and are inside of D'. Let x_0 be the center of D'. Then

$$\max_{t_i \in A} d(t_i, x_0) < \min_{t_j \in B} d(t_j, x_0)$$

It follows from Lemma 3.1 that $A \prec B$ does not hold. Hence $C(T, e) \neq \emptyset$.

Case (b). If A_e consists of only one point t_i , then C(T, e) contains a point (ε, t_i) for sufficiently small $\varepsilon > 0$. If B_e consists of only one point, then $C(T, e^*) \neq \emptyset$.

Case (c). The sets A_e , B_e do not give the natural partition of $\{t_1, \dots, t_{d+2}\}$. So we have $\overline{A}_e \cap \overline{B}_e = \emptyset$ by the uniqueness of the natural partition. We can find a (d-1)-dimensional sphere D' such that $D' \supset B_e$ and all the points of A_e are inside of D'. Then $C(T, e) \neq \emptyset$, since

$$\max_{t_i\in A_e} d(t_i, x_0) < \min_{t_j\in B_e} d(t_j, x_0)$$

for the center x_0 of D'.

EXAMPLE 4.2. Let $T_0 = (t_1, \dots, t_{d+1}) \in (\mathbf{R}^d)^{d+1}$ and $t_{d+2} = 0$. We choose and fix T_0 in such a way that $T = (t_1, \dots, t_{d+1}, t_{d+2})$ satisfies the assumption in Lemma 4.1. It follows from Lemmas 2.5 and 4.1 that $S(T, e) \neq \emptyset$ for every $e \in \mathscr{E}_{d+1}$. Let μ be a measure on $E = \mathbf{R}_+ \times \mathbf{R}^d$ satisfying (1.3) such that $\mu(S(T_0, e)) > 0$ for every $e \in \mathscr{E}_{d+1}$. Let us define $\tilde{\mu}$ in the following way. We make $\tilde{\mu} = \mu$ on $E \setminus \bigcup_{i=1}^{d+1} S_{t_i}$. First notice that μ satisfies the consistency condition (2.11) for n = d + 1. Using the notations in the proof of Lemma 2.2, let A be the matrix $M_{d+1}(2^{d+1} - 1)$ and b be the vector in (2.11). Let c be the $(2^{d+1} - 2)$ -vector every component of which is $\mu(\bigcap_{i=1}^{d+1} S_{t_i})$. Choose $\varepsilon \neq 0$ such that every component of the solution \mathbf{x} of

$$A\mathbf{x} = \mathbb{D} - (1 + \varepsilon)\mathbb{C}$$

is positive. It suffices to make $|\varepsilon|$ small enough. Now let

$$\tilde{\mu} \Big(igcap_{i=1}^{d+1} S_{\iota_i} \Big) = (1+arepsilon) \mu \Big(igcap_{i=1}^{d+1} S_{\iota_i} \Big)$$

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and let $\tilde{\mu}(S(T_0, e))$ for $e \in \mathscr{E}_{d+1} \setminus \{(1, \dots, 1)\}$ be given by the solution x. There exists a measure $\tilde{\mu}$ with these $\tilde{\mu}(S(T_0, e))$, $e \in \mathscr{E}_{d+1}$. We have $\tilde{\mu}(S_t) < \infty$ for all $t \in \mathbb{R}^d$. Let $\{X(t)\}$ and $\{\tilde{X}(t)\}$ be the Chentsov type $S \alpha S$ random fields associated with μ and $\tilde{\mu}$, respectively. From the construction

$$ilde{\mu}(ilde{S}(T_{\scriptscriptstyle 0},e))=\mu(ilde{S}(T_{\scriptscriptstyle 0},e)) \qquad ext{for all } e \in \mathscr{E}_{d+1}ackslash\{(1,\,\cdots,\,1)\}\,.$$

It follows that $(X(t_1), \dots, X(t_{d+1}))$ and $(\tilde{X}(t_1), \dots, \tilde{X}(t_{d+1}))$ have different distributions but they have common *d*-dimensional marginal distributions.

EXAMPLE 4.3. An interesting problem is whether there are two measures μ and $\tilde{\mu}$ satisfying (1.3) such that the Chentsov type $S\alpha S$ random fields $\{X(t)\}$ and $\{\tilde{X}(t)\}$ on \mathbb{R}^d associated with μ and $\tilde{\mu}$, respectively, have identical *d*-dimensional distributions but different (d + 1)-dimensional distributions. We do not know the answer to this problem for general *d* yet. But, in case d = 1, we can construct such measures.

Let $E = \mathbf{R}_+ \times \mathbf{R}^i$. Let μ be such that $\mu(S_t) = \mu(S_{-t}) < \infty$ and $\mu(S_t)$ is a continuous increasing function of t > 0. Suppose, further, that μ is mutually absolutely continuous with the Lebesgue measure. Let $\tilde{\mu}$ be a measure concentrated on $\mathbf{R}_+ \times \{0\}$ such that

$$\tilde{\mu}(S_t) = \tilde{\mu}(S_t \cap (\mathbf{R}_+ \times \{0\})) = \mu(S_t) \,.$$

Then $\{X(t)\}$ and $\{\tilde{X}(t)\}$ have common 1-dimensional distributions. Let $0 < t_1 < t_2$. Then $\mu(S_{t_1} \cap S_{t_2}^c) > 0$ but $\tilde{\mu}(S_{t_1} \cap S_{t_2}^c) = 0$, which implies that $(X(t_1), X(t_2))$ and $(\tilde{X}(t_1), \tilde{X}(t_2))$ have different distributions.

Our technique in this paper works in finding determinism of random fields on \mathbf{R}^{d} of a similar sort.

THEOREM 4.4. Let μ be a measure on $\mathbf{R}_+ \times \mathbf{R}^d$ satisfying $\mu(C_t) < \infty$ for every $t \in \mathbf{R}^d$ and let $Y(\cdot)$ be the $S \alpha S$ random measure associated with μ . Let

$$X(t) = Y(C_t)$$
 for $t \in \mathbf{R}^d$.

Then, for any n > d, any n-dimensional distribution of $\{X(t)\}$ is determined by its d-dimensional marginal distributions.

Proof. The non-degenerate case is dealt with Proposition 3.6 and Lemma 2.2. The degenerate case is obvious. \Box

Finally we remark that, if μ is invariant under translation in \mathbf{R}^{d} ,

then $\{X(t); t \in \mathbb{R}^d\}$ in Theorem 4.4 is a homogeneous random field constructed geometrically.

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