

## Local Theorems for Euclidean Motions. I

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### 1. Introduction

Let  $\mathbb{G}$  be some topological group,  $\nu$  a probability measure on  $\mathbb{G}$ ,  $\nu^{*n}$  the  $n$ -fold convolution of  $\nu$ . The value  $\nu^{*n}(\mathcal{D})$  is interesting for some applications where  $\mathcal{D}$  is any bounded domain in  $\mathbb{G}$ . Every estimate of the value  $\nu^{*n}(\mathcal{D})$ ,  $n \rightarrow \infty$  will be called a local theorem (l.t.) on the group  $\mathbb{G}$ . Accordingly the estimate of the value  $\nu^{*n}(\mathcal{D}_1)/\nu^{*n}(\mathcal{D}_2)$  for any bounded domains  $\mathcal{D}_1$  and  $\mathcal{D}_2$  is called a relation local theorem (r.l.t.). Local theorems have a short history. They have been proved in the case of additive groups  $\mathbb{R}^1$  and  $\mathbb{R}^d$  by Shepp [11] and Stone [12] respectively in 1965. Kazdan [5] was the first who began to study r.l.t. for the group of Euclidean plane motions. This result refers to a particular case of a measure  $\nu$  whose support contains four motions  $\mathcal{A}$ ,  $\mathcal{A}^{-1}$ ,  $\mathcal{B}$ ,  $\mathcal{B}^{-1}$ . This result was obtained by functional methods and didn't possess a probability character. A r.l.t. has also been proved for local compact unimodal groups by Le Pege [6] and for a wide class of commutative groups by Stone [13].

The first l.t. for the group of Euclidean plane motions was reported in the 3rd Soviet Japanese symposium on probability in 1975 [10]. Independently P. Baldi, Ph. Bougerol, P. Crépel obtained analogous results under wider conditions for Euclidean motions of the  $d$ -dimensional space  $\mathbb{R}^d$  using the Plancherel formula for the Euclidean group of motions [1, 2].

By  $\mathbb{M}(d)$  we denote the group of Euclidean motions  $\mathbb{R}^d$ .

Some details from [10] were given by the author in [8] and [9], especially with reference to the group  $\mathbb{M}(2)$ . In particular the concept of "lattice motions" of  $\mathbb{M}(2)$  in [9] together with l.t. [1] or with the theorem of the present work give l.t. for any case of measures  $\nu$  on  $\mathbb{M}(2)$  with finite support. It is clear that Kazdan's theorem follows from the general cases of l.t. with a finite support of the measure  $\nu$ . The purpose of the present work is to obtain the same [1] by the direct probability method. This method uses a variant of Stone's l.t. [12] applied to independent and nonidentically distributed values and it is sensitive to complementary conditions. In particular, the method can give some other variants of l.t. that do not follow from [1]. They will be obtained in the second part of this work.

## 2. The Formulation of the Main Result

The group of Euclidean motions  $\mathbf{M}(d)$  is the cross-group of the additive group  $\mathbb{R}^d$  with the rotations group  $\mathbf{SO}(d)$ . Every element of this group is uniquely determined by the pair  $(x, \alpha)$ , where  $x$  is the shift and  $\alpha$  is the rotation. The product of the motions is defined by  $(x, \alpha)(y, \beta) = (x + y\alpha, \alpha\beta)$ , where  $\alpha\beta$  is the product of the rotations and  $y\alpha$  is the image  $y \in \mathbb{R}^d$  under  $\alpha \in \mathbf{SO}(d)$ . Let  $\nu$  be a measure defined on  $\mathbf{M}(d)$ . Taking its natural desintegration we obtain for every  $\alpha \in \mathbf{SO}(d)$  the measure  $\nu(\alpha)$  on  $\mathbb{R}^d$  and the measure  $\mathbf{v}$  on  $\mathbf{SO}(d)$  which is the conditional distribution  $\nu$ . The random elements  $\xi_1, \xi_2, \dots, \xi_n, \dots$  in  $\mathbf{M}(d)$  correspond to the random independent pairs  $(x_1, \alpha_1), \dots, (x_n, \alpha_n), \dots$ . For the pairs  $(x_k, \alpha_k)$ ,  $k=1, 2, \dots$  the random values  $x_k$  and  $\alpha_k$  may be dependent. Then the product  $\xi_1 \xi_2 \dots \xi_n$  corresponds to  $(x_1 + x_2\alpha_1 + \dots + x_n\alpha_1 \dots \alpha_{n-1}, \alpha_1 \dots \alpha_n)$ . Consequently, the product  $\xi_1 \xi_2 \dots \xi_n$  may be reduced to the study of the joint distribution of  $\mathcal{S}_n = x_1 + x_2\alpha_1 + \dots + x_n\alpha_1 \dots \alpha_{n-1} \in \mathbb{R}^d$  and  $\alpha_1 \alpha_2 \dots \alpha_n \in \mathbf{SO}(d)$ . Under general conditions the product  $\alpha_1 \alpha_2 \dots \alpha_n$  converges to the uniform distribution on  $\mathbf{SO}(d)$  and becomes asymptotically independent with  $\mathcal{S}_n$ . Consequently the l.t. in  $\mathbf{M}(d)$  is equivalent to the estimate of the probability  $\mathbf{P}\{\mathcal{S}_n \in \mathcal{D}\}$ . Therefore we will estimate this probability. Let  $\xi$  be arbitrary  $\mathbf{M}(d)$  - valued random variable having the distribution  $\nu$  determined by the pair  $(x, \alpha)$ . By  $x(\alpha)$  we denote the conditional variable with fixed  $\alpha$ ,  $\mathbf{v}$  - the distribution  $\alpha$ ,  $\mathcal{F}_\alpha$  - the distribution  $x$ ,  $\nu(\alpha)$  - the distribution  $x(\alpha)$  and  $\nu'(\alpha)$  - the distribution obtained by symmetrization of  $\nu(\alpha)$ .

**Theorem** (l.t.) *Let  $\xi_1, \xi_2, \dots, \xi_n, \dots$  be independent  $\mathbf{M}(d)$  - valued random variables with identical distribution  $\nu$  and satisfying the conditions:*

1. *the  $n$ -fold convolution  $\nu^{*n}$  converges weakly as  $n \rightarrow \infty$  to the uniform distribution on  $\mathbf{SO}(d)$ .*
2. *The distribution  $\mathcal{F}_\alpha$  on  $\mathbb{R}^d$  has finite second moments.*
3. *There exists some number  $n_0 > 0$  and a set  $K$ ,  $K \subseteq \mathbf{SO}(d)$  of measure  $\nu(K) > 0$  so that for any  $\alpha \in K$  the support of distribution  $\nu^{*n_0}(\alpha)$  does not contain any hyperplanes. In what follows for the sake simplicity we shall suppose  $n_0 = 1$ .*

*Then for any bounded Jordan measurable domain  $\mathcal{D} \in \mathbb{R}^d$ ,*

$$\mathbf{P}\{\mathcal{S}_n \in \mathcal{D}\} = c \text{mes}(\mathcal{D}) n^{-d/2} + O(n^{-d/2}) \quad (1)$$

*where  $\text{mes}(\mathcal{D})$  is the Lebesgue measure of the domain  $\mathcal{D}$ ,  $c$  is a constant which does not depend on  $\mathcal{D}$  and will be evaluated at the end of this paper.*

It is easy to see that the formulations of the proposition [1] and our theorem coincide.

The main tool of the direct probability method for studying local properties of the sum  $\mathcal{S}_n$  is the reduction to the conditional sums  $\mathcal{S}_n(\alpha_1, \dots, \alpha_n)$  of independent but not identically distributed random variables where

$$\mathcal{S}_n(\alpha_1, \dots, \alpha_n) = x(\alpha_1) + x(\alpha_1)\alpha_1 + \dots + x(\alpha_n)\alpha_1 \dots \alpha_{n-1} \quad (2)$$

for fixed collections  $(\alpha_1, \alpha_2, \dots, \alpha_n)$ .

For fixed  $(\alpha_1, \alpha_2, \dots, \alpha_n)$  the conditional variables are independent but not identically distributed. If for the "majority" of these sums the l.t. in  $\mathbb{R}^d$  may be

applied, i.e. the value  $\mathbf{P}\{\mathcal{S}_n(\alpha_1, \dots, \alpha_n) \in \mathcal{D}\}$  may be effectively estimated, then we have by integration

$$\mathbf{P}\{\mathcal{S}_n \in \mathcal{D}\} = \int \mathbf{P}\{\mathcal{S}_n(\alpha_1, \dots, \alpha_n) \in \mathcal{D}\} \mathbf{v}(d\alpha_1) \dots \mathbf{v}(d\alpha_n). \quad (3)$$

In spite of the essential difference of the methods that have been applied here and in [1] they use effectively the same condition 3. For example whenever this condition or analogous ones are absent the l.t. has not yet been proved, if  $d \geq 3$  and the measure  $\mathbf{v}$  has a finite support without additional conditions.

Finally we remark that by virtue of condition 2 we get the following condition 2': the variables  $z(\alpha)$ ,  $\alpha \in \mathbf{SO}(d)$  have second moments which are absolutely integrable. The expectation  $\mathbf{M}z(\alpha) = m(\alpha)$  has second moments also. We denote the characteristic function of  $z(\alpha)$  by  $f_\alpha(t)$ .

### 3. A Special Formulation of the l.t. in $\mathbb{R}^d$

To apply the l.t. to the sum (2) we need the following variant of Stone's l.t.

**Lemma 1.** *Let  $\xi_1(\tau), \xi_2(\tau), \dots, \xi_n(\tau), \dots$  be a sequence of  $\mathbb{R}^d$ -valued nonidentically distributed variables depending on the parameter  $\tau$  on a measurable compact  $\mathbb{T}$  with probability measure  $\mu$ . For any  $\tau \in \mathbb{T}$  the variables  $\xi_k(\tau)$  are independent and possess second moments with respect to the measure  $\mu$ . By  $f_n(t, \tau_n)$  we denote the characteristic function of the sum  $\xi_1(\tau_n) + \dots + \xi_n(\tau_n)$ . If for all  $n$ ,  $n > 0$  there exists a set  $\mathcal{K}_n \subseteq \mathbb{T}$ ,  $\mu(\mathcal{K}_n) \rightarrow 1$  and we have uniformly for any sequence  $\{\tau_n\}$ ,  $\tau_n \in \mathcal{K}_n$ :*

1) *the distribution of the variable  $\frac{1}{\sqrt{n}}\{\xi_1(\tau_n) + \dots + \xi_n(\tau_n) - m(n, \tau)\}$  converges*

*weakly as  $n \rightarrow \infty$  to the normal distribution with mean zero and covariance matrix  $\Omega$  where*

$$m(n, \tau_n) = M\{\xi_1(\tau_n) + \dots + \xi_n(\tau_n)\},$$

2) *the integral  $\sqrt{n} \int_b^B |f_n(t, \tau_n)| dt$  tends to zero as  $n \rightarrow \infty$  for any positive numbers  $b, B$ ,*

3) *the integral  $\sqrt{n} \int_{\frac{A}{\sqrt{n}} \leq \|t\| \leq b} |f_n(t, \tau_n)| dt$  tends to zero as  $n \rightarrow \infty$ ,  $A \rightarrow \infty$ ,  $b \rightarrow 0$ ,*

*then for any bounded Jordan measurable domain  $\mathcal{D}$  and for all sequences  $\{\tau_n\}$ ,  $\tau_n \in \mathcal{K}_n$  we have*

$$\begin{aligned} \mathbf{P}\{\xi_1(\tau_n) + \dots + \xi_n(\tau_n) \in \mathcal{D}\} &= (\sqrt{n})^{-d} \int_{\mathcal{D}} \mathbf{N}_\Omega\left(\frac{y}{\sqrt{n}} - \frac{m(n, \tau_n)}{\sqrt{n}}\right) dy \\ &+ o_n(A, b, \tau_n)(\sqrt{n})^{-d}, \end{aligned} \quad (4)$$

where  $\mathbf{N}_\Omega(x)$  is the density of the normal distribution with mean zero and covariance matrix  $\Omega$ , and the value  $o_n(A, b, \tau_n)$  tends to zero uniformly as  $n \rightarrow \infty$ ,  $A \rightarrow \infty$ ,  $b \rightarrow 0$ .

The scheme of Stone's proof [12] remains valid when applied to the proof of Lemma 1. Since we intend to apply Lemma 1 to the sum  $\mathcal{L}_n(\alpha_1, \dots, \alpha_n)$  it is necessary to represent the terms of the sum in accordance with Lemma 1. In order to do this we take the compact space  $\mathbb{T} = \mathbb{S}\mathbb{O}(d)^\infty$  with measure  $\mu$ , where  $\mathbb{S}\mathbb{O}(d)^\infty$  is equal to a countable product of the group  $\mathbb{S}\mathbb{O}(d)$  and the measure  $\mu$  is the direct product of the measure  $\mathbf{v}$  on  $\mathbb{S}\mathbb{O}(d)$ . For any  $\tau = (\alpha_1, \alpha_2, \dots) \in \mathbb{T}$  we construct the sequence of  $\mathbb{R}^d$ -valued variables

$$\xi_1(\tau) = \varkappa(\alpha_1), \quad \xi_2(\tau) = \varkappa(\alpha_2)\alpha_1, \dots, \xi_n(\tau) = \varkappa(\alpha_n)\alpha_1 \dots \alpha_{n-1} \quad (5)$$

which are independent and nonidentically distributed.

Then for any  $\tau \in \mathbb{T}$  we can define the random variables

$$\mathcal{L}_n(\tau) = \mathcal{L}_n(\alpha_1, \dots, \alpha_n) = \xi_1(\tau) + \dots + \xi_n(\tau).$$

We shall apply the equality (4) to the sum  $\mathcal{L}_n(\alpha_1, \dots, \alpha_n)$ . For this it is enough to verify that the conditions 1), 2), 3) are valid for  $\mathcal{L}_n(\alpha_1, \dots, \alpha_n)$ .

#### 4. The Normal Convergence of the Sequence $\frac{1}{\sqrt{n}} \{\mathcal{L}_n(\tau_n) - \mathbf{M}(\mathcal{L}_n(\tau_n))\}$

Condition 1 will hold for  $\mathcal{L}_n(\alpha_1, \dots, \alpha_n)$  if the distribution of the sequence  $\frac{1}{\sqrt{n}} \{\mathcal{L}_n(\tau_n) - \mathbf{M}(\mathcal{L}_n(\tau_n))\}$ ,  $\tau_n \in \mathcal{K}_n$  converges to the normal distribution with density  $N_{\Omega_1}(x)$  where  $\mathcal{K}_n$  are some subsets of  $\mathbb{S}\mathbb{O}(d)^\infty$ ,  $\mu(\mathcal{K}_n) \rightarrow 1$ . We need the following auxilliary proposition:

**Lemma 2.** *Let  $\alpha_1, \alpha_2, \dots$  be independent  $\mathbb{S}\mathbb{O}(d)$ -valued variables with the same distribution  $\mathbf{v}$ . Let  $f(x)$  be the linear sum of elements taken from the matrix of nontrivial irreducible representation of  $\mathbb{S}\mathbb{O}(d)$ ,  $a(x)$  a real function on  $\mathbb{S}\mathbb{O}(d)$  with  $\int |a(x)| d\mathbf{v} \leq \infty$ . Then*

$$\frac{1}{n} \sum_1^n a(\alpha_k) f(\alpha_1 \dots \alpha_{k-1}) \quad (6)$$

tends to zero weakly.

*Proof.* At first we suppose that  $a(x)$  has second moments and we shall prove that

$$\frac{1}{n^2} \mathbf{M} \left\{ \sum_1^n a(\alpha_k) f(\alpha_1 \dots \alpha_{k-1}) \right\}^2 \rightarrow 0. \quad (7)$$

The sum (7) is equal to

$$\begin{aligned} & \frac{1}{n^2} \sum_1^n \mathbf{M} \alpha^2(\alpha_k) \mathbf{M} f^2(\alpha_1 \dots \alpha_{k-1}) \\ & + \frac{1}{n^2} \sum_{1 \leq k < m \leq n} \mathbf{M} \alpha(\alpha_k) a(\alpha_m) f(\alpha_1 \dots \alpha_{k-1}) f(\alpha_1 \dots \alpha_{k-1} \alpha_k \dots \alpha_{m-1}). \end{aligned}$$

Since the function  $f(\alpha)$  is continuous,  $|f(\alpha)| < c_1$  on  $\mathbb{S}\mathbb{O}(d)$ . Let us set  $\mathbf{M} a^2(\alpha) = c_2$ ,  $\mathbf{M} |a(\alpha)| = c_3$ . Under condition 1) for the measure  $\nu$  it was proved in [7] that  $|\mathbf{M} f(\alpha_1 \dots \alpha_k)| < c_4 \lambda^k$  where  $c_4, \lambda$  are constants not depending on  $\alpha$ , and in addition  $0 < \lambda < 1$ . Therefore the first sum remains less than  $\frac{1}{n} c_2 c_1^2$  and the second one is less than  $\frac{1}{n^2} \sum_{1 \leq k < m \leq n} c_3^2 c_1 c_4 \lambda^{m-k}$ . Consequently (7) is true, and  $\frac{1}{n} \sum a(\alpha_k) f(\alpha_1 \dots \alpha_{k-1})$  tends to zero weakly if  $\mathbf{M} a^2(\alpha) < \infty$ . If  $\mathbf{M} a^2(\alpha) = \infty$  we can write  $a(\alpha) = a_N(\alpha) + \varepsilon_N(\alpha)$  where  $|a_N(\alpha)| < N$  and  $\mathbf{M} |\varepsilon_N| < \delta$ . Then the sum (6) is equal to

$$\frac{1}{n} \sum_1^n a_N(\alpha_k) f(\alpha_1 \dots \alpha_{k-1}) + \frac{1}{n} \sum_1^n \varepsilon_N(\alpha_k) f(\alpha_1 \dots \alpha_{k-1}).$$

It is proved above that the first sum tends to zero and the second sum is less than  $c_1 \frac{1}{n} \sum_1^n |\varepsilon_N(\alpha_k)|$ . But  $c_1 \frac{1}{n} \sum_1^n |\varepsilon_N(\alpha_k)|$  tends to  $c_1 \mathbf{M} |\varepsilon_N|$  which is less than  $c_1 \delta$ . Since  $\delta$  may be chosen arbitrarily we get the required result.

**Lemma 3.** *Under the conditions of Lemma 2 the value  $\frac{1}{n} \sum_1^n a(\alpha_k) f^2(\alpha_1 \dots \alpha_{k-1})$  converges weakly to  $\mathbf{M} a(\alpha) \int f^2(\alpha) d\alpha$  where the integral  $\int f^2(\alpha) d\alpha$  is taken with respect to the invariant measure on  $\mathbb{S}\mathbb{O}(d)$ .*

*Proof.* The function  $f^2(\alpha)$  is an element of the tensor square of the corresponding representation. Therefore

$$f^2(\alpha) = \gamma_0 + \sum_1^d \gamma_k f_k(\alpha) \quad (8)$$

where  $\gamma_1$  is constant,  $f_k(\alpha)$  are elements of the some irreducible nontrivial representations  $\mathbb{S}\mathbb{O}(d)$  and  $\int f^2(\alpha) d\alpha = \gamma_0$  because  $\int f_k(\alpha) d\alpha = 0$ . Consequently applying Lemma 2 to the sum (8) we obtain Lemma 3.

By  $(z_1(\alpha), \dots, z_n(\alpha))$  we denote the coordinates of the random vector  $z(\alpha)$  with mean zero, and by  $g(w) = \|g_{ij}(w)\|$  we denote the orthogonal matrix of the rotation  $w \in \mathbb{S}\mathbb{O}(d)$ . The map  $w \rightarrow g(w)$  is continuous and  $g(w, w_2) = g(w_1) g(w_2)$ ,  $g(w)$  is an irreducible representation of  $\mathbb{S}\mathbb{O}(d)$ . Let  $\mathbf{M}(\alpha)$  be the covariance matrix of  $z(\alpha)$  and  $z(\alpha)w$  be the image of  $z(\alpha) \in \mathbb{R}^d$  under  $w \in \mathbb{S}\mathbb{O}(d)$ . Then the covariance of  $z(\alpha)w$  is equal to  $g(w)^T \mathbf{M}(\alpha) g(w)$ , where  $g(w)^T$  is the transposed matrix. However it follows that the covariance matrix for the random variable

$\frac{1}{\sqrt{n}} \{\mathcal{L}_n - \mathbf{M} \mathcal{L}_n(\tau)\}$ , where  $\tau = (\alpha_1, \alpha_2, \dots)$  is equal to

$$\frac{1}{n} \sum_{k=1}^n g^T(\alpha_1 \dots \alpha_{k-1}) \mathbf{M}(\alpha_k) g(\alpha_1 \dots \alpha_{k-1}) = \Omega(n; \alpha_1, \dots, \alpha_n). \quad (9)$$

By virtue of condition 2' the expectation  $\mathbf{M} \{\mathbf{M}(\alpha)\} = \|\sigma_{ij}\|$  exists. Set  $\sigma = \frac{1}{d} (\sigma_{11} + \dots + \sigma_{dd})$ . It follows from condition 1, that  $\sigma > 0$ .

**Lemma 4.** *The value  $\Omega(n; \alpha_1, \dots, \alpha_n)$  converges weakly to  $\sigma \mathbb{I}$  where  $\mathbb{I}$  is the identical matrix.*

*Proof.* The coordinates of the vector  $\varkappa(\alpha) w$  are equal to

$$\varkappa_j(\alpha, w) = \varkappa_1(\alpha) g_{j1}(w) + \dots + \varkappa_d(\alpha) g_{jd}(w), \quad j=1, 2, \dots, d.$$

Then the  $(i, j)$ -th element of the matrix  $\Omega$  is equal to

$$\frac{1}{n} \sum_1^n \left\{ \sum_{m,r} M_{\varkappa_m(\alpha_k) \varkappa_r(\alpha_k)} g_{mi}(\alpha_1 \dots \alpha_{k-1}) g_{rj}(\alpha_1 \dots \alpha_{k-1}) \right\}. \quad (10)$$

The functions  $g_{ij}(w)$  are the elements of an irreducible non-trivial unitary representation for which the following orthogonal relations are well-known:

$$\int g_{ij}^2(w) dw = \frac{1}{d}; \quad \int g_{mi}(w) g_{rj}(w) = 0 \quad \text{if } (m, i) \neq (r, j).$$

The functions  $g_{ij}(w) g_{rj}(w)$  are the elements of the tensor square representation of  $g(w)$ . Therefore to every sum of the kind

$$\frac{1}{n} \sum_{k=1}^n M_{\varkappa_m(\alpha_k) \varkappa_r(\alpha_k)} g_{mi}(\alpha_1 \dots \alpha_{k-1}) g_{rj}(\alpha_1 \dots \alpha_{k-1})$$

Lemmas 2 or 3 may be applied. This proves Lemma 4.

Now we will indicate the set  $\mathcal{K}_n \subset \mathbb{SO}(d)^\infty$ ,  $\mu(\mathcal{K}_n) \rightarrow 1$ , for which the distribution of the variable  $\frac{1}{\sqrt{n}} \{ \mathcal{S}_n(\tau_n) - \mathbf{M} \mathcal{S}_n(\tau_n) \}$ ,  $\tau_n \in \mathcal{K}_n$  converges weakly and uniformly to the normal distribution  $\mathbf{N}_{\Omega_1}(x)$ ,  $\Omega_1 = \sigma \mathbb{I}$ .

At first we set  $\varphi_\lambda(x) = \sup_{w, \alpha_0} \int_{\|x - \alpha_0\| > 2} \|x w\|^2 dv(\alpha)$  where  $\|x\| = \sqrt{x_1^2 + \dots + x_d^2}$ .

Evidently the function is measurable and  $\varphi_\lambda(x) \leq \int \|x\|^2 dv(\alpha)$ .

Since the right hand side of the inequality is integrable in  $\mathbf{v}$ , we have for all  $\lambda > 0$ ,  $\int \varphi_\lambda(x) dv = \varphi(\lambda) < c$  where the constant  $c$  does not depend on  $\lambda$ . For any  $\alpha$  the function  $\varphi_\lambda(\alpha)$  decreases to zero as  $\lambda \rightarrow \infty$ . Consequently  $\varphi(\lambda) \rightarrow 0$  as  $\lambda \rightarrow \infty$ ,

and there exists a sequence  $\lambda_N$ ,  $\varphi(\lambda_N) = \varepsilon_N \downarrow 0$ . Since the sum  $\frac{1}{n} \sum_1^n \varphi_\lambda(\alpha) \rightarrow \varphi(\lambda)$  in  $\mu$  measure, we may define the sets  $\mathcal{K}_N$ ,  $\mathcal{K}_N \subset \mathcal{K}_{N+1} \subset \mathbb{SO}(d)^\infty$ ,  $\mu(\mathcal{K}_N) \uparrow 1$  such that for any  $(\alpha_1, \alpha_2, \dots) \in \mathcal{K}_N$  the following inequality holds:  $\limsup \frac{1}{n} \sum_{k=1}^n \varphi_{\lambda_N}(\alpha_k) \leq \varepsilon_{N-1}$ .

In virtue of the decreasing  $\varphi_\lambda(\alpha)$  as  $\lambda \rightarrow \infty$  for any  $\alpha$  and  $\tau = (\alpha_1, \alpha_2, \dots) \in \mathcal{K}_N$  we have  $\limsup \frac{1}{n} \sum_{k=1}^n \varphi_{\ln n}(\alpha_k) < \varepsilon_{N+m}$  for any integer  $N$ ,  $m > 0$  or  $\limsup \frac{1}{n} \sum_{k=1}^n \varphi_{\ln n}(\alpha_k) = 0$ . This implies that

$$\lim \frac{1}{n} \sum_1^n \int_{\|x - m(\alpha_k) \alpha_1 \dots \alpha_{k-1}\| > \varepsilon \sqrt{n}} \|x\|^2 dv(\alpha_k) \alpha_1 \dots \alpha_{k-1} \leq \limsup \frac{1}{n} \sum_1^n \varphi_{\ln n}(\alpha_k) = 0. \quad (11)$$

The last inequality holds for arbitrary but fixed  $\mathcal{K}_N$ , any  $\varepsilon > 0$  and from some positive integer  $N_\varepsilon$  on.

The right-hand side of the inequality (11) is a necessary and sufficient condition for the convergence in the weak sense of  $\frac{1}{\sqrt{n}} \{ \mathcal{S}_n(\tau_n) - \mathbf{M} \mathcal{S}_n(\tau_n) \}$  to the normal distribution under  $\tau_n \in \mathcal{K}_N$ .

On the other hand the covariance matrix from Lemma 4  $\frac{1}{\sqrt{n}} \{ \mathcal{S}_n(\tau_n) - \mathbf{M} \mathcal{S}_n(\tau_n) \}$  converges to  $\Omega_1 = \sigma \Pi$ . Since the normal distribution with zero is determined by the covariance matrix uniquely and  $\mu(\mathcal{K}_N) \rightarrow 1$  as  $N \rightarrow \infty$ , condition 1) of Lemma 1 for the variables (5) is satisfied.

*Remark.* It is well-known by [4] that the distribution  $\frac{1}{\sqrt{n}} \mathcal{S}_n$  converges weakly to the normal one with density  $\mathbf{N}_\Omega(x)$ . The matrix  $\Omega$  was defined in [4], too. On other hand we have

$$\mathcal{S}_n = \frac{1}{\sqrt{n}} \{ \mathcal{S}_n(\tau) - \mathbf{M} \mathcal{S}_n(\tau) \} + \frac{1}{\sqrt{n}} \mathbf{M} \mathcal{S}_n(\tau)$$

for any collection  $(\alpha_1, \alpha_2, \dots) = \tau$ .

It was shown above that the distribution  $\frac{1}{\sqrt{n}} \{ \mathcal{S}_n(\tau_n) - \mathbf{M} \mathcal{S}_n(\tau_n) \}$ ,  $\tau_n \in \mathcal{K}_n$  converges to the normal one with density  $\mathbf{N}_{\Omega_1}(x)$ . The distribution of the variable  $\frac{1}{\sqrt{n}} \mathbf{M} \mathcal{S}_n(\tau) = \frac{1}{\sqrt{n}} \sum m(\alpha_k) \alpha_1 \dots \alpha_{k-1}$  converges to the normal distribution with the density  $\mathbf{N}_{\Omega_2}(x)$  [4] since by condition 2', the value of  $\int \|m(\alpha)\|^2 d\mathbf{v}$  is bounded, and the random variables  $\alpha_1, \alpha_2, \dots$  are independent with common distribution  $\mathbf{v}$ , satisfying condition 1. Consequently  $\mathbf{N}_\Omega(x) = \mathbf{N}_{\Omega_1 + \Omega_2}(x)$ . This equality will be used to calculate the constant  $c$  in the theorem.

### 5. Proof of Conditions 2), 3) for the Sum $\mathcal{S}_r(\alpha_1, \dots, \alpha_n)$

To check conditions 2), 3) for the sum  $\mathcal{S}_n(\alpha_1, \dots, \alpha_n)$  we need the following proposition.

**Lemma 5.** *Let  $f_\alpha(t)$  be the characteristic function of the variable  $x(\alpha)$ . Put  $w_k = (\alpha_1 \dots \alpha_k)^{-1}$ ,  $w_0$  - the identity of  $\mathbf{SO}(d)$ . Then for any segment  $[b, B]$ ,  $0 < b < B$  there exists  $\theta$ ,  $0 < \theta < 1$  such that the inequality*

$$\mathbf{M} \left\{ \prod_{k=1}^n |f_{\alpha_k}(t w_{k-1})| \right\} < \theta^n$$

holds uniformly for any  $t \in [b, B]$  as  $n \rightarrow \infty$ .

*Proof.* For  $\alpha \in K$  (see condition 3.) the support of the distribution  $\nu(\alpha)$  does not contain any hyperplane. Thus, if  $\alpha \in K$ , then the function  $|f_\alpha(t)|$  may be equal to 1

only for the points laying on some lines  $t_\alpha$  passing through the origin  $\mathbb{O}$ . Let  $\beta$  be a random variable on  $\mathbb{SO}(d)$  with distribution  $\mathbf{v}$ . The convergence of  $\mathbf{v}^{*n}$  to the uniform distribution on  $\mathbb{SO}(d)$  involves the convergence of  $(\mathbf{v}w)^{*n}$  to one under fixed  $w \in \mathbb{SO}(d)$ . Then  $\mathbf{M}_\beta |f_\alpha(t\beta w)| < 1$ , since the distribution  $v(\alpha)$  must be degenerate in the contrary case. Thus  $\mathbf{M}_\beta |f_\alpha(t\beta w)| = v(t, \alpha, w) < 1$  for  $t \neq 0$ ,  $\alpha \in K$ ,  $v(t, \alpha, w)$  is a continuous function and  $\sup_w v(t, \alpha, w) = v(t, \alpha) < 1$ ,  $\mathbf{M}_\alpha v(t, \alpha) = v(t) < 1$  for  $t \neq 0$ ,  $\alpha \in K$ .

Hence we can find a  $\theta_1$ ,  $0 < \theta_1 < 1$  such that  $v(t) < \theta_1$  for all  $t \in [b, B]$ . Since  $|f_\alpha(t)| \leq 1$  we have the following inequalities:

$$\mathbf{M} \left\{ \prod_{k=1}^n |f_{\alpha_k}(t w_{k-1})| \right\} \leq \mathbf{M} \left\{ \prod_{k=1}^{\lfloor n/2 \rfloor} |f_{\alpha_{2k}}(t w_{2k-1})| \right\} \quad (12)$$

and

$$\begin{aligned} & \mathbf{M}_{\alpha_{2k}} (\mathbf{M}_{2k-1} |f_{\alpha_{2k}}(t \alpha_{2k-1}^{-1} \cdot \alpha_{2k-2}^{-1} \dots \alpha_1^{-1})|) \\ & = \mathbf{M}_{\alpha_{2k}} v(t, \alpha_{2k}, w_{2k-2}) \leq \mathbf{M}_{\alpha_{2k}} v(t, \alpha_{2k}) = v(t) < \theta_1. \end{aligned} \quad (13)$$

We remark that the random pairs  $(\alpha_{2k}, \alpha_{2k-1})$  and  $(\alpha_{2m}, \alpha_{2m-1})$  are independent for any integers  $k, m$ ,  $k \neq m$ . Thus it follows from (13) and  $\mathbf{M}_\alpha \mathbf{M}_\beta |f_\alpha(t\beta w)| < \theta_1$  that the right hand side of (12) can be estimated recursively, i.e.

$$\mathbf{M} \left\{ \prod_{k=1}^{\lfloor n/2 \rfloor} |f_{\alpha_{2k}}(t w_{2k-1})| \right\} \leq \mathbf{M} \left\{ \prod_{k=1}^{\lfloor n/2 \rfloor - 1} |f_{\alpha_{2k}}(t w_{2k-1})| \right\} \theta_1 \leq \dots \leq \theta_1^{\lfloor n/2 \rfloor}.$$

This completes the proof.

To verify that condition 2) is fulfilled for the sum  $\mathcal{L}_n(\alpha_1, \dots, \alpha_n)$  it suffices to show that

$$\mathbf{M} \left\{ \sqrt{n} \int_{b \leq \|t\| \leq B} \prod |f_{\alpha_k}(t w_{k-1})| dt \right\} \rightarrow 0.$$

But by Lemma 5 this value is less than  $\theta^n \sqrt{n}(B-b)$ ,  $0 < \theta < 1$  and consequently it tends to zero. To verify that condition 3) is fulfilled it is enough to show that

$$\sqrt{n} \mathbf{M} \left\{ \int_{\frac{A}{n} \leq \|t\| \leq b} \prod |f_{\alpha_k}(t w_{k-1})| dt \right\}$$

tends to zero as  $n, A \rightarrow \infty$ ,  $b \rightarrow 0$ . For this we put  $\mathbf{M} |f_\alpha(t)| = u(t)$ ,  $u(t w_i) = \max_w u(t w)$ . We know that the variables  $\alpha_k$ ,  $k=1, 2, \dots$  are independent and distributed identically, hence we have

$$\mathbf{M} \prod_{k=1}^n |f_{\alpha_k}(t w_{k-1})| \leq \left\{ \mathbf{M} \prod_{k=1}^{n-1} |f_{\alpha_k}(t w_{k-1})| \right\} u(t w_n) \leq u^n(t w_n). \quad (14)$$

It is regrettable that the function  $u(t)$  cannot be a characteristic function. However,  $\{\mathbf{M} |f_\alpha(t)|\}^2 = u(t)^2 \leq \mathbf{M} |f_\alpha(t)|^2 = \mathbf{U}(t)$  where  $\mathbf{U}(t)$  is the characteristic function of the distribution  $\int v(\alpha) d\mathbf{v}$ . This distribution has finite second moments and by condition 3) does not contain any hyperplanes. Hence it follows



that for some neighbourhood of zero we have the estimate

$$U(t) < 1 - c \|t\|^2.$$

Since  $\|t\| = \|t w\|$  for any rotation  $w \in \mathbb{SO}(d)$ , we also have

$$U(t w) < 1 - c \|t\|^2. \tag{15}$$

Consequently it follows from (14) and (15) that

$$\sqrt{n} \mathbf{M} \int_{\frac{A}{\sqrt{n}} \leq \|t\| \leq b} \prod |f_x(t w_{k-1})| dt \leq \sqrt{n} \int_{\frac{A}{\sqrt{n}} \leq \|t\| \leq b} U^{[n/2]}(t) dt. \tag{16}$$

According to (15) standard tools can be applied to estimate the integral of the power of the characteristic function in the segment  $[A/\sqrt{n}, b]$ . In that case we obtain  $\sqrt{n} \int U^{[n/2]}(t) dt \rightarrow 0$  as  $n, A \rightarrow \infty, b \rightarrow 0$ .

Consequently, this completes the proof of condition 3) for the sum  $\mathcal{S}_n(\alpha_1, \dots, \alpha_n)$ .

**6. Esseen's Inequality for the Sum  $\mathcal{S}_n(\alpha_1, \dots, \alpha_n)$ .**

Let  $\xi_1, \xi_2, \dots, \xi_n, \dots$  be a sequence of random variables which are independent and identically distributed in  $\mathbb{R}^d$ . In the paper by Esseen [3] the following inequality was obtained for the sum  $\mathcal{S}_n = \xi_1 + \dots + \xi_n$

$$\sup_{\alpha \in \mathbb{R}^d} \mathbf{P} \{ \mathcal{S}_n \in \Sigma_r + \alpha \} = c(d, r) \left\{ \sup_{u \geq r} u^{-2} \sum_1^n \chi_k(u) \right\}^{-d/2} \tag{17}$$

where  $\Sigma_r$  is the sphere of radius  $r$  with the centre at the origin;  $c(d, r)$  is a constant depending only on  $d$  and  $r$ ,  $\mathbf{P}_k$  is the distribution of  $\xi_k$ ,  $\chi_k(u) = \inf_{\|t\|=1} \int_{\|x\| < u} (t, x)^2 d\mathbf{P}'_k$  where the distribution  $\mathbf{P}'_k$  is obtained by symmetrization of  $\mathbf{P}_k$ . Let us apply the inequality (17) to the sum

$$\mathcal{S}_n(w; \alpha_1, \dots, \alpha_n) = x_1(\alpha_1) w_1 + \dots + x_n(\alpha_n) w_n$$

where the rotations  $w_1, \dots, w_n$  are arbitrary but fixed. Since for any rotation  $w \in \mathbb{SO}(d)$  the equality  $(v(\alpha) w)' = v'(\alpha) w$  holds, we shall have

$$\begin{aligned} \inf_{\|t\|=1} \int_{\|x\| \leq \lambda_0} (t, x)^2 d(v(\alpha) \tau)' &= \inf_{\|t\|=1} \int_{\|x\| \leq \lambda_0} (t \tau^{-1}, x)^2 d v'(\alpha) \\ &= \inf_{\|t\|=1} \int_{\|x\| \leq \lambda_0} (t, x)^2 d v'(\alpha). \end{aligned} \tag{18}$$

Thus Esseen's inequality for the sum  $\mathcal{S}_n(w; \alpha_1, \dots, \alpha_n)$  does not depend on the arbitrary rotations  $w_1, \dots, w_n \in \mathbb{SO}(d)$ . The function  $\psi(\alpha, \lambda) = \inf_{\|t\|=1} \int_{\|x\| \leq \lambda} (t, x)^2 d v'(\alpha)$  is measurable with respect to  $\alpha$  and is decreasing with respect to  $\lambda$  for every fixed  $\alpha$ . By condition 3. if  $\alpha \in K$  then  $\liminf_{\lambda \rightarrow \infty} \psi(\alpha, \lambda) > 0$ .

Consequently we can find a subset  $K_0 \subset K$  with positive measure  $\mathbf{v}(K_0) = p_0 > 0$  and some  $\lambda_0$  such that simultaneously for all  $\alpha \in K_0$ ,  $\lambda \geq 0$  the inequality

$$\inf_{\|t\|=1} \int_{\|x\|<\lambda} (t, x)^2 d\nu'(x) \geq \delta_0 > 0 \quad (19)$$

is fulfilled.

**Lemma 6.** *Let us suppose that among  $\alpha_1, \dots, \alpha_n$  the value  $\alpha_k \in K_0$  is taken on  $p_0 n/2$  times at least. Then there exists a sphere  $\Sigma_r$  and a constant  $c(p_0, \delta, d, r)$  such that*

$$\sup_{z \in \mathbb{R}^d} \mathbf{P} \{ \mathcal{L}_n(\alpha_1, \dots, \alpha_n) \in \Sigma_r + z \} < c(p_0, \delta, d, r) n^{-d/2}.$$

In fact this can be obtained by employing (18) and (19) in the inequality (17) for the sum  $\mathcal{L}_n(\alpha_1, \dots, \alpha_n)$ .

The next lemma follows from the preceding one.

**Lemma 7.** *Under the conditions of Lemma 6 for any bounded domain  $\mathcal{D}$  in  $\mathbb{R}^d$  there exists a constant  $c(\mathcal{D}, p_0, \delta, d)$  such that*

$$\sup_{z \in \mathbb{R}^d} \mathbf{P} \{ \mathcal{L}_n(\alpha_1, \dots, \alpha_n) \in \mathcal{D} + z \} \leq c(\cdot) n^{-d/2}. \quad (20)$$

Furthermore the measure of all trajectories  $(\alpha_1, \dots, \alpha_n)$  for which the inequality (20) may not be true is less than  $\exp\{-c(p_0)n\}$  and the constant  $c(p_0) > 0$  depends only on  $p_0$ .

*Proof.* In fact any bounded domain  $\mathcal{D}$  of diameter  $L$  may be covered by the finite number of spheres of the radius  $r$ . This number depends only on  $L$  and  $r$ . Consequently we can apply Lemma 6. The second part of Lemma 7 follows by elementary combinatorial estimations.

## 7. The Proof of the Theorem

It follows from above that we can define some sequence of sets  $\mathcal{K}_n \subset \mathbb{S}\mathbb{O}(d)^\infty$ ,  $\mu(\mathcal{K}_n) \rightarrow 1$  as  $n \rightarrow \infty$  such that the conditions of Lemma 1 are fulfilled for the sum  $\mathcal{L}_n(\alpha_1, \dots, \alpha_n)$  if  $\tau_n = (\alpha_1, \dots, \alpha_n, \dots) \in \mathcal{K}_n$ . Then for every number  $n$  we shall represent the compact  $\mathbb{S}\mathbb{O}(d)^\infty$  as the sum of the sets  $\mathcal{K}_n$ ,  $\mathcal{L}_n$  and  $\mathcal{E}_n$  where  $\mathcal{L}_n$  is a set of  $\tau$  such that  $\tau = (\alpha_1, \dots, \alpha_n, \dots) \in \mathbb{S}\mathbb{O}(d)^\infty$ , so that the value  $\alpha_k \in K_0$  is taken on less than  $p_0 n/2$  times among the  $n$  first components,  $\mathcal{E}_n$  is the rest of the elements  $\mathbb{S}\mathbb{O}(d)$ . It is obvious that  $\mu(\mathcal{E}_n) \leq \varepsilon_n$ . By virtue of Lemma 7 we have  $\mu(\mathcal{L}_n) < \exp(-c(p_0)n)$ . It is clear from the foregoing definition of the measure  $\mu$  on  $\mathbb{S}\mathbb{O}(d)^\infty$  that integral (3) can be written

$$\mathbf{P} \{ \mathcal{L}_n \in \mathcal{D} \} = \int \mathbf{P} \{ \mathcal{L}_n(\tau) \in \mathcal{D} \} d\mu = \int_{\mathcal{K}_n} + \int_{\mathcal{L}_n} + \int_{\mathcal{E}_n}. \quad (21)$$

By Lemma 6 and the evident inequality  $0 \leq \mathbf{P} \{ \mathcal{L}_n(\tau) \in \mathcal{D} \} \leq 1$  we have

$$\int_{\mathcal{K}_n} = (\sqrt{n})^{-d} \int_{\mathcal{K}_n} \left( \int_{\mathcal{D}} \mathbf{N}_{\Omega_1}(\cdot) dy \right) d\mu + \left( \int \theta_n(\cdot) d\mu \right) n^{-d/2}.$$

Since  $\theta_n$  tends to zero uniformly on  $\mathcal{K}_n$  as  $n \rightarrow \infty$  the second integral is less than  $\varepsilon'_n n^{-d/2}$  where  $\varepsilon'_n \rightarrow 0$ .

It is obvious that

$$\int_{\mathcal{K}_n} = n^{-d/2} \mathbf{M}_\tau(\cdot) - n^{-d/2} \int_{\mathcal{L}_{n+\varepsilon_n}}, \quad \text{where} \quad \int_{\mathcal{L}_{n+\varepsilon_n}} \mathbf{N}_{\Omega_1}(0) \text{ mes } \mathcal{D} \cdot \varepsilon''_n, \varepsilon'''_n \rightarrow 0.$$

To evaluate  $M_\tau(\cdot)$  it suffices to change the integration order and to note that  $\text{mes}(\mathcal{D})$  and  $\mathbf{N}_{\Omega_1}(x)$  are bounded. Therefore

$$\mathbf{M}_\tau(\cdot) = \text{mes } \mathcal{D} \cdot \left\{ \lim_n \mathbf{M}_\tau \left[ \mathbf{N}_{\Omega_1} \left( \frac{\mathbf{M}_{\mathcal{L}_n(\tau)}}{\sqrt{n}} \right) \right] \right\} + \varepsilon''_n, \quad \varepsilon''_n \rightarrow 0.$$

Thus we obtain

$$\mathbf{P}\{\mathcal{L}_n \in \mathcal{D}\} = \text{mes } \mathcal{D} \left\{ \lim_n \mathbf{M}_\tau[\cdot] \right\} n^{-d/2} + o_n(1) n^{-d/2}.$$

According to the remark above the distribution of the variable  $\mathbf{M} \frac{\mathcal{L}_n(\tau)}{\sqrt{n}}$  converges uniformly to the normal distribution with the density  $\mathbf{N}_\Omega(x)$ .

Consequently the constant  $\lim_n \mathbf{M}_\tau \left[ \mathbf{N}_{\Omega_1} \left( \frac{\mathbf{M}_{\mathcal{L}_n(\tau_n)}}{\sqrt{n}} \right) \right] = c$  is equal to

$$\int_{\mathbb{R}^d} \mathbf{N}_{\Omega_1}(x) \mathbf{N}_{\Omega_2}(x) dx = \mathbf{N}_\Omega(0), \text{ where } \mathbf{N}_\Omega(x) \text{ is the density of the limit distribution}$$

of  $\frac{1}{\sqrt{n}} \mathcal{L}_n$ . Thus the theorem is proved.

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