

On Separation between Metric Observers in Segal's Compact Cosmos

Alexander Levichev¹, Andrey Palyanov²

¹Sobolev Institute of Mathematics SB RAS, Novosibirsk, Russia

²A. P. Ershov Institute of Informatics Systems SB RAS, Novosibirsk, Russia

Email: alevichev@gmail.com

Received 14 July 2015; accepted 14 November 2015; published 17 November 2015

Copyright © 2015 by authors and Scientific Research Publishing Inc.

This work is licensed under the Creative Commons Attribution International License (CC BY).

<http://creativecommons.org/licenses/by/4.0/>



Open Access

Abstract

A certain class K of GR homogeneous spacetimes is considered. For each pair E, \tilde{E} of spacetimes from K , $\tilde{E} = g(E)$ where conformal transformation g is from $G = SU(2,2)$. Each E (being $U(2)$ or its double cover, as a manifold) is interpreted as related to an observer in Segal's universal cosmos. The definition of separation d between E and \tilde{E} is based on the integration of the conformal factor of the transformation g . The integration is carried out separately over each region where the conformal factor is no less than 1 (or no greater than 1). Certain properties of $d = d(E, \tilde{E})$ are proven; examples are considered; and possible directions of further research are indicated.

Keywords

Separation between Spacetimes, Segal's Universal Cosmos, Conformal Group Action on $U(2)$, DLF-Theory

1. Motivation and Introduction

The first author has been interested in GR ("GR" is for General Relativity) research for quite a while and he concentrated on a few most symmetric spacetimes ([1], [2], and more). Later (see [3], [4]) he has become a strong believer in Segal's Chronometric Theory (see [5], electronic archive arranged by Levichev), and he is attempting to modify Segal's Theory (see [6], a key publication). The collaboration of the two current authors is based on their mutual interest in Penrose-Hameroff approach to consciousness (see its update in [7], [8]). Specifically, we are putting forward an alternative definition of separation between space-times. In [9], the original definition was based on bringing up a Newtonian limit in GR. Our definition has been introduced in

[10], [11], and we now present it in much more detail.

Recall the Lie group $U(2)$ as the totality of all two-by-two matrices z (with complex entries allowed) satisfying

$$z^* z = \mathbf{1}, \tag{1.1}$$

where z^* is the transpose and complex conjugate of z , and $\mathbf{1}$ is the unit matrix. Now, define the Lie group $G = SU(2, 2)$ as consisting of all four-by-four matrices g (with complex entries allowed) satisfying

$$g^* S g = S, \tag{1.2}$$

where S is the diagonal matrix $\text{diag}\{1, 1, -1, -1\}$. Recall the well-known linear-fractional G -action on $U(2)$:

$$g(z) = (Az + B)(Cz + D)^{-1}, \tag{1.3}$$

where a matrix g from G is determined by its 2×2 blocks A, B, C, D .

In Table I of [12], the matrices L_{ij} are chosen as basic vectors of the (fifteen-dimensional) Lie algebra $su(2, 2)$, whereas L_{ij} are the corresponding vector fields on $U(2)$. The vector fields L_{ij} are determined by the G -action (1.3). As explained in [12], subscripts i, j take on $-1, 0, 1, 2, 3, 4$, and the convention $L_{ji} = -L_{ij}$ (resulting in $L_{ji} = -L_{ij}$) holds.

The Lorentzian inner product on $U(2)$ is introduced in such a way that left-invariant vector fields $X_0 = L_{-10}, X_1 = L_{14} - L_{23}, X_2 = L_{24} - L_{31}, X_3 = L_{34} - L_{12}$ form an orthonormal basis (following [12], we use $+, -, -, -$ signature). The resulting product on $U(2)$ is bi-invariant (see [6]), and $\langle \mathbf{a}, \mathbf{b} \rangle$ below denotes the Lorentzian inner product of tangent vectors \mathbf{a}, \mathbf{b} at a point z of $U(2)$. The spacetime thus obtained is denoted by E_0 (the meaning of the subscript will become clear in the next section). Transformations (1.3) are conformal in E_0 (a word of caution: this spacetime has been denoted as D in [6]). As it follows from Table I of [12], the vector fields $L_{-10}, L_{12}, L_{13}, L_{14}, L_{23}, L_{24}, L_{34}$, generate isometries in E_0 . The corresponding subgroup K in G consists of all matrices (1.2) with $B = C = 0$.

For what follows, it is instrumental to introduce a certain bi-invariant Riemannian inner product on $U(2)$. To do so, we recall that vector fields X_0, X_1, X_2, X_3 constitute a basis of the Lie algebra $u(2)$. This algebra is a direct sum of its center with $su(2)$. Namely, X_0 generates the center, whereas X_1, X_2, X_3 are basic vectors in $su(2)$. The Riemannian metric is determined by the demand on left-invariant vector fields X_0, X_1, X_2, X_3 to form an orthonormal basis. The corresponding Riemannian space is denoted E_R . Again, E_R is $U(2)$, as a manifold. Our (\mathbf{a}, \mathbf{b}) below denotes the Riemannian inner product of tangent vectors \mathbf{a}, \mathbf{b} at a point z of $U(2)$. In the forthcoming sections the corresponding volume form on E_R will be instrumental. From Table I of [12], it follows that the group K acts as a group of both Lorentzian and Riemannian isometries.

Notice that, as a group, $U(2)$ is not a direct product of its center with the subgroup $SU(2)$. The double cover $E^{(2)}$ of $U(2)$ is the direct product $S^1 \times S^3$ (with S^3 represented by $SU(2)$). The covering map sends (e^{it}, u) into the matrix $e^{it}u$ in $U(2)$. The corresponding Lorentzian metric on $E^{(2)}$ is of the form

$$(dt)^2 - (du)^2. \tag{1.4}$$

Here the variable t is along S^1 whereas $(du)^2$ is for the standard Riemannian metric on S^3 . More details are given in our Appendix A, where u denotes a matrix from $SU(2)$. Our Appendix B is dedicated to a certain one-parameter group of transformations (1.3).

It is well-known ([12], [13]) that the (above introduced) covering map is a Lorentzian isometry. Infinitesimal G -action on $E^{(2)}$ is presented in Table I of [12]. It is known (see [13]) that action (1.3) can be lifted to a global conformal G -action on $E^{(2)}$. Using the corresponding commutative diagram, one can show that the lifted action of the group K is as follows:

$$(e^{it}, u) \rightarrow ((\det A)e^{it}, (\det A)^{-1} AuD^{-1}). \tag{1.5}$$

Also, it is easily verifiable that for the Riemannian metric

$$(dt)^2 + (du)^2 \tag{1.6}$$

on $S^1 \times S^3$ the (above specified) covering map is a Riemannian isometry from $E^{(2)}$ onto E_R .

It makes sense to mention how a suitable version of the Einstein static universe, E_{uc} , can be introduced in the

context of our work (the subscript uc is for universal cover). To be more precise, E_{uc} should be called *universal cosmos* ([14]) or *Segal's universal cosmos* ([13]). The universal cover E_{uc} of E_0 is $R^1 \times S^3$, topologically. The G -action on E_0 is canonically lifted to the G_{uc} -action on E_{uc} (the latter action preserves the causal structure of E_{uc}). In a cosmological model based on E_{uc} , there is a conformal invariant R , interpreted as the radius of a three-dimensional (physical) space S^3 . I. Segal (in [14] and in other publications) has put this R for the (long wanted by Dirac and others) third fundamental constant additionally to the speed of light and to the Planck's constant. It is known that to model particles on E_{uc} , one can start with the world E_0 , a compact one. The respective property is called automatic periodicity ([15], p. 202), and it allows us to only deal with compact spacetimes E_0 and $E^{(2)}$ (which explains our *compact Segal's cosmos* terminology).

More precisely, we deal with two classes of spacetimes: \mathbf{K}_1 and \mathbf{K}_2 . Roughly speaking, the first class is obtained by application of all transformations (3) to E_0 ; details follow in the next section (where \mathbf{K}_2 will also be defined).

2. On the Notion of Separation between Spacetimes: The Main Definition and Related Properties

The separation (or *distance*) $d(x, y)$ will be defined for any pair x, y of spacetimes from \mathbf{K}_1 (or from \mathbf{K}_2).

As mentioned, the totality of all isometries in E_0 is the group K of all matrices (1.2) with $B = C = 0$. Each member E of the class \mathbf{K}_1 will be now put in correspondence with an element x of the homogeneous space G/K . Namely, each element (or *coset*) x of G/K is specified by an element g from G : $x = gK = \{gk : k \in K\}$. One and the same x can be determined by another element (say, g_1) from G : $gK = g_1K$. For such a pair (g, g_1) , there exists such k from K , that $g_1 = gk$. When the subgroup K is viewed as an element of G/K , denote K as x_0 . This x_0 we put into correspondence with E_0 (which has been described in Section 1). As a manifold, each element x of \mathbf{K}_1 is $U(2)$. In what follows, we use $\langle \mathbf{a}, \mathbf{b} \rangle$ (rather than $\langle \mathbf{a}, \mathbf{b} \rangle_0$) to denote the (Lorentzian) inner product of vectors \mathbf{a}, \mathbf{b} from the tangent space $T(E_0)$ at z . This inner product has been introduced in our Section 1. To define spacetime E corresponding to a coset $x = gK$, it is enough to specify the inner product $\langle \cdot, \cdot \rangle_E$, see (2.2) below. Such a transformation g is conformal in E_0 . Namely, given vectors \mathbf{a}, \mathbf{b} from the tangent space $T(E_0)$ at z , the inner product $\langle g_*\mathbf{a}, g_*\mathbf{b} \rangle$ at $g(z)$ of their images (under the tangent map g_*) satisfies

$$\langle g_*\mathbf{a}, g_*\mathbf{b} \rangle = h(z)\langle \mathbf{a}, \mathbf{b} \rangle. \tag{2.1}$$

The everywhere positive function $h(z)$ is known as the *square of conformal coefficient*. Frequently, we will simply refer to this $h(z)$ as to a *conformal coefficient*. Given vectors \mathbf{a}, \mathbf{b} from the tangent space $T(E)$ at z , their inner product *can be defined* as follows:

$$\langle \mathbf{a}, \mathbf{b} \rangle_E = [h(g^{-1}(z))]^{-1} \langle \mathbf{a}, \mathbf{b} \rangle, \tag{2.2}$$

where $\langle \mathbf{a}, \mathbf{b} \rangle$ is calculated at z . Notice, that h here is the (above mentioned) function on $U(2)$ determined by g . It is easy to show that (2.2) is equivalent to the condition for g to be an isometry between $(E_0, \langle \cdot, \cdot \rangle)$ and $(E, \langle \cdot, \cdot \rangle_E)$:

$$\langle g_*\mathbf{a}, g_*\mathbf{b} \rangle_E = \langle \mathbf{a}, \mathbf{b} \rangle. \tag{2.3L}$$

Here the right hand side of (2.3L) is calculated in E_0 at z , and it defines the inner product $\langle g_*\mathbf{a}, g_*\mathbf{b} \rangle_E$ of vectors $g_*\mathbf{a}$ and $g_*\mathbf{b}$ at $g(z)$. To avoid verification of (2.2)-(2.3L) equivalence, we define the Lorentzian metric in E in terms of (2.3L). Similarly, we define the following Riemannian metric in E :

$$(g_*\mathbf{a}, g_*\mathbf{b})_E = (\mathbf{a}, \mathbf{b}). \tag{2.3R}$$

where the positive definite inner product in the right hand side of (2.3R) has been introduced in our Section 1.

Let us show that, given a coset x in G/K , (2.3L) correctly defines a Lorentzian metric on $U(2)$, whereas (2.3R) correctly defines a Riemannian metric on $U(2)$:

Scholium 2.1. The inner product (2.3L) (respectively, the inner product (2.3R)) is independent of the choice of g which represents a coset x .

Proof. If x is represented as g_1K , then $g_1 = gk$ where k is a certain element of the group K . Given such a representation, the analogue of (2.3R) is

$$\{g_{1*}\mathbf{a}, g_{1*}\mathbf{b}\}_E = \langle \mathbf{a}, \mathbf{b} \rangle \quad (2.4)$$

where the right hand side of (2.4) is calculated in E_0 at z , and it defines the inner product $\{g_{1*}\mathbf{a}, g_{1*}\mathbf{b}\}_E$ of vectors $g_{1*}\mathbf{a}$ and $g_{1*}\mathbf{b}$ at $g_1(z)$. We have to show that (2.4) introduces the same metric structure on $U(2)$ as (2.3L) does. To do so, we rewrite (2.3L) in the form of

$$\{(gk)_*\mathbf{a}, (gk)_*\mathbf{b}\}_E = \langle k_*\mathbf{a}, k_*\mathbf{b} \rangle \quad (2.5)$$

where the right hand side is calculated at $k(z)$, and the left hand side is calculated at $g(k(z))$.

However, $\langle k_*\mathbf{a}, k_*\mathbf{b} \rangle$ in (2.5) equals $\langle \mathbf{a}, \mathbf{b} \rangle$ in (2.4) since k is a Lorentzian isometry in E_0 . Comparison of (2.4) with (2.5) finishes the proof. The verification process, that (2.3R) is independent of representative, copies the one for (2.3L). \square

Let us notice (see [16]) that each $(E, \langle \cdot, \cdot \rangle_E)$ can be interpreted as a spacetime corresponding to a certain (global) observer. A word of caution: [16] treats the universal cover of E_0 whereas we only deal with compact spacetimes here.

Remark 2.2. We have thus defined the class \mathbf{K}_1 of spacetimes. Our class \mathbf{K}_2 can be similarly introduced in terms of the Lorentzian manifold $E^{(2)}$ (the 2-cover of E_0) with the (lifted) G -action on $E^{(2)}$.

Given a (1.3)-transformation g of (Lorentzian) E_0 , define the following subsets of E_0 :

$$T_g^+ = \{z : h_g(z) \geq 1\}, \quad (2.6)$$

$$T_g^- = \{z : h_g(z) \leq 1\}; \quad (2.7)$$

here $h_g(z)$ is the square of the conformal coefficient at z of the transformation g . A (non-negative) number d_g is defined as follows:

$$d_g = \ln \left\{ \frac{[\int h_g^2]}{[V(T_g^+)]} \right\} + \ln \left\{ \frac{[\int h_g^{-2}]}{[V(T_g^-)]} \right\}, \quad (2.8)$$

where $V(S)$ is for the volume of a set S in E_R (with the volume form introduced in our Section 1). Clearly, expressions inside the logarithms in (2.8) can be interpreted as corresponding cumulative distortions of the original metric structure in E_0 . To be sure of convergence of all of the integrals involved, it is enough to mention that each of the two integrands is a continuous function over the corresponding region of integration, whereas each of the regions (2.6), (2.7) is a compact set.

To further deal with (2.8), we now proceed with more technicalities. Clearly, d_g can be viewed as $\ln \{(ac)/(bd)\}$ with $b = V(T_g^+)$, $d = V(T_g^-)$. The integration in a is over T_g^+ , whereas in c the integration is over T_g^- .

Examples of integrals a, b, c, d evaluations are given in our Appendix C (see Theorem C.5). Notice that

$$d_g = 0 \text{ if and only if } g \text{ is an isometry of } E_0, \quad (2.9)$$

which follows from (2.8) because in this case $h_g(z) = 1$, a constant function on $U(2)$. As a result of $h_g(z) = 1$, each of the two terms in the sum (2.8) is zero.

Scholium 2.3. Given the (1.3)-transformation g and isometries k_1, k_2 , the following holds:

$$d_w = d_g, \quad (2.10)$$

where $w = k_1^{-1}gk_2$.

Proof. To prove (2.10), we will now show that each of the four numbers (a, b, c , and d) remain the same when we switch from d_g to d_w . Namely:

$$\begin{aligned} T_w^+ &= \{z : h_w(z) \geq 1\} = k_2^{-1} \left(\{k_2(z) : h_m(k_2(z)) \geq 1\} \right) \\ &= k_2^{-1} \left(\{k_2(z) : h_g(k_2(z)) \geq 1\} \right), \text{ where } m = k_1^{-1}g; \text{ hence } T_w^+ = k_2^{-1}(T_g^+). \end{aligned}$$

Similarly, $T_w^- = k_2^{-1}(T_g^-)$. Hence, $b = V(T_g^+) = V(T_w^+)$, and $d = V(T_g^-) = V(T_w^-)$, due to the K-invariance of the volume form.

Let us now use the variable $k_2(z)$ in the integral a of $(h_w)^2$ over T_w^+ : the integrand is then $(h_m(k_2(z)))^2 = (h_g(k_2(z)))^2$, the region of integration is $k_2(T_w^+) = T_g^+$, and there is no extra factor in the integrand since k_2 is a transformation from the group K. Similarly, number c remains the same when we switch from d_g to d_w . \square

Now, if two cosets are represented as $x = g_1K$, $y = g_2K$, define the distance $d(x, y)$ as

$$d(x, y) = d_g, \quad (2.11)$$

where $g = (g_1)^{-1}g_2$. The number $d(x, y)$ is independent of representatives since if x is represented by g_1k_1 , and y is represented by g_2k_2 , then for $w = (g_1k_1)^{-1}g_2k_2 = k_1^{-1}gk_2$, $d_w = d_g$ according to (2.10).

Corollary 2.4. In the above settings, $d(x, y) = d(x_0, q)$ where $x_0 = K$ and $q = g_1^{-1}g_2K$.

A word of caution: we use the term *distance* but we are not sure that the corresponding triangle inequality holds (even locally) for (2.11). However, we prove (below) that (2.11) is symmetric: $d(x, y) = d(y, x)$, and G -invariant:

$$d(f(x), f(y)) = d(x, y), \quad (2.12)$$

for arbitrary f from G (where we have in mind the canonical action of G in G/K).

As regards G -invariance, one can think of a possible relation of our definition (2.11) to the canonical inner product in the symmetric space G/K . This we do not discuss here.

Scholium 2.5. The distance (2.11) is symmetric: $d(x, y) = d(y, x)$.

Proof. As justified by our Corollary 2.4, assume that $x = K$ and $y = gK$. Define

$$T_m^+ = \{\tilde{z} : h_m(\tilde{z}) \geq 1\}, T_m^- = \{\tilde{z} : h_m(\tilde{z}) \leq 1\},$$

where $m = g^{-1}$, and where we use \tilde{z} (rather than z , as before) to denote a matrix in $U(2)$. Tilde (below) indicates that computations are performed in E rather than in E_0 .

For $d_m = \ln \left\{ \frac{\int (h_m)^2}{\int \tilde{V}(T_m^+)} \right\} + \ln \left\{ \frac{\int (h_m)^{-2}}{\int \tilde{V}(T_m^-)} \right\}$, the following is true:

$\tilde{b} = \tilde{V}(T_m^+) = \tilde{V}(g(T_g^-)) = V(T_g^-) = d$ since, due to (2.3R), g is an isometry between the two Riemannian spaces.

Similarly, $\tilde{d} = \tilde{V}(T_m^-) = \tilde{V}(g(T_g^+)) = V(T_g^+) = b$. The new integrand is then $(h_g(z))^{-2}$, the new region of integration is T_g^- , and there is no extra factor in the integrand since g is an isometry between the two Riemannian spaces in question. We have thus proven that $\tilde{a} = c$. Similarly, $\tilde{c} = a$. We have thus proven the equality $(\tilde{a}/\tilde{b})(\tilde{c}/\tilde{d}) = (a/b)(c/d)$, which results in $d(x, y) = d(y, x)$, the symmetry property of the distance between spacetimes. \square

3. Concluding Remarks and Future Research Insights

Examples of integrals a, b, c, d evaluations (in case of a certain one-parameter group of conformal transformations) are given in our Appendix C. It is of interest to know whether Theorem C.5 holds for other transformations from $G = SU(2, 2)$. Evaluations in Appendix C indicate that definition (2.11) of distance between spacetimes seems to be quite a working one. As part of future research, it will be of interest to apply our definition in the case where the original spacetime is F (here we refer to the DLF-theory, [6]). In that case, the underlying manifold is (non-compact!) $U(1, 1)$, rather than $U(2)$. Preliminary calculations indicate that a

conformal coefficient might be unbounded. We will thus have to deal with improper 4D integrals, and the question of convergence will have to be studied first.

References

- [1] Guts, A.K. and Levichev, A.V. (1984) On the Foundations of Relativity Theory. *Doklady Akademii Nauk SSSR*, **277**, 253-257. (In Russian)
- [2] Levichev, A.V. (1989) On the Causal Structure of Homogeneous Lorentzian Manifolds. *General Relativity & Gravity*, **21**, 1027-1045. <http://dx.doi.org/10.1007/BF00774087>
- [3] Levichev, A.V. (1993) The Chronometric Theory by I. Segal Is the Crowning Accomplishment of Special Relativity. *Izvestiya Vysshikh Uchebnykh Zavedenii Fizika*, **8**, 84-89. (In Russian)
- [4] Levichev, A.V. (1995) Mathematical Foundations and Physical Applications of Chronometry. In: Hilgert, J., Hofmann, K. and Lawson, J., Eds., *Semigroups in Algebra, Geometry, and Analysis*, de Gruyter Expositions in Mathematics, Berlin, 77-103. <http://dx.doi.org/10.1515/9783110885583.77>
- [5] <http://dedekind.mit.edu/segal-archive/index.php>
- [6] Levichev, A.V. (2011) Pseudo-Hermitian Realization of the Minkowski World through DLF Theory. *Physica Scripta*, **83**, 1-9. <http://dx.doi.org/10.1088/0031-8949/83/01/015101>
- [7] Hameroff, S. and Penrose, R. (2014) Consciousness in the Universe: A Review of the “Orch OR” Theory. *Physics of Life Reviews*, **11**, 39-78. <http://dx.doi.org/10.1016/j.plrev.2013.08.002>
- [8] Hameroff, S. and Penrose, R. (2014) Reply to Criticism of the “Orch OR Qubit”—“Orchestrated Objective Reduction” Is Scientifically Justified. *Physics of Life Reviews*, **11**, 94-100. <http://dx.doi.org/10.1016/j.plrev.2013.11.013>
- [9] Penrose, R. (1992) Gravity and Quantum Mechanics. In: Gleiser, R.J., Kozameh, C.N. and Moreschi, O.M., Eds., *General Relativity and Gravitation 13. Part 1: Plenary Lectures 1992. Proceedings of the 13th International Conference on General Relativity and Gravitation*, Cordoba, 28 June-4 July 1992, 179-189.
- [10] Levichev, A. and Palyanov, A. (2014) On a Modification of the Theoretical Basis of the Penrose-Hameroff Model of Consciousness. In: *International Conference MM-HPC-BBB-2014*, Abstracts, Sobolev Institute of Mathematics SB RAS, Institute of Cytology and Genetics, Siberian Branch of the Russian Academy of Sciences, Novosibirsk, 49.
- [11] Levichev, V. and Palyanov, A.Yu. (2014) On a Notion of Separation between Space-Times. In: *Geometry Days in Novosibirsk. Abstracts of the International Conference*, Sobolev Institute of Mathematics, Siberian Branch of the Russian Academy of Sciences, Novosibirsk, 110.
- [12] Segal, I.E., Hans, P.J., Bent, Ø., Paneitz, S.M. and Speh, B. (1981) Covariant Chronogeometry and Extreme Distances: Elementary Particles. *PNAS*, **78**, 5261-5265. <http://dx.doi.org/10.1073/pnas.78.9.5261>
- [13] Werth, J.-E. (1986) Conformal Group Actions and Segal’s Cosmology. *Reports on Mathematical Physics*, **23**, 257-268. [http://dx.doi.org/10.1016/0034-4877\(86\)90023-6](http://dx.doi.org/10.1016/0034-4877(86)90023-6)
- [14] Segal, I.E. (1976) *Mathematical Cosmology and Extragalactic Astronomy*. Academic Press, New York.
- [15] Branson, T.P. (1987) Group Representations Arising from Lorentz Conformal Geometry. *Journal of Functional Analysis*, **74**, 199-291. [http://dx.doi.org/10.1016/0022-1236\(87\)90025-5](http://dx.doi.org/10.1016/0022-1236(87)90025-5)
- [16] Segal, I.E. (1984) Evolution of the Inertial Frame of the Universe. *Nuovo Cimento*, **79B**, 187-191. <http://dx.doi.org/10.1007/BF02748970>
- [17] Kon, M. and Levichev, A. (2015) Towards Analysis in Space-Time Bundles Based on Pseudo-Hermitian Realization of the Minkowski Space. In Preparation.

Appendix A: Parameterizations of $U(2)$ and $E^{(2)}$

The following presentation for $E^{(2)}$, the 2-cover of $U(2)$, has been widely used in the literature. Consider the direct sum $E^6 = E^2 \oplus E^4$ of two Euclidean spaces: E^2 with rectangular coordinates u_{-1}, u_0 , and E^4 with rectangular coordinates u_1, u_2, u_3, u_4 . Each “event” in $E^{(2)}$ is a 6-tuple $(u_{-1}, u_0, u_1, u_2, u_3, u_4)$, satisfying

$$u_{-1}^2 + u_0^2 = 1, \tag{A1}$$

and

$$u_1^2 + u_2^2 + u_3^2 + u_4^2 = 1. \tag{A2}$$

Clearly, $E^{(2)}$ is $S^1 \times S^3$, topologically. The earlier introduced e^{it} (see Section 2) is $u_{-1} + iu_0$, whereas the matrix u from $SU(2)$ is specified as follows:

$$u = \begin{bmatrix} u_4 + iu_3 & u_2 + iu_1 \\ iu_1 - u_2 & u_4 - iu_3 \end{bmatrix}. \tag{A3}$$

The covering map from $E^{(2)}$ onto $U(2)$ takes the pair $(u_{-1} + iu_0, u)$ into the matrix $(u_{-1} + iu_0)u$, an element z of the group $U(2)$:

$$z = (u_{-1} + iu_0)u. \tag{A4}$$

Given a matrix z in $U(2)$, the factors $(u_{-1} + iu_0)$ and u are defined up to a sign, only. In terms of E^6 , it is helpful to consider a pseudo-Euclidean metric

$$(du_{-1})^2 + (du_0)^2 - (du_1)^2 - (du_2)^2 - (du_3)^2 - (du_4)^2, \tag{A5L}$$

and an Euclidean metric

$$(du_{-1})^2 + (du_0)^2 + (du_1)^2 + (du_2)^2 + (du_3)^2 + (du_4)^2. \tag{A5R}$$

It is known (see [14], p. 40) that the restriction of (A5L) onto $E^{(2)} = S^1 \times S^3$ coincides with metric (1.4) of our Section 1. Similarly, the restriction of (A5R) onto $E^{(2)}$ coincides with metric (1.5).

Appendix B: The Case of a Certain One-Parameter Group of Conformal Transformations

This group consists of all (1.3)-transformations g of the form:

$$g = \begin{bmatrix} A & B \\ C & D \end{bmatrix}, \tag{B.1}$$

with

$$A = \begin{bmatrix} 1 & 0 \\ 0 & c \end{bmatrix}, B = \begin{bmatrix} 0 & 0 \\ s & 0 \end{bmatrix}, C = \begin{bmatrix} 0 & s \\ 0 & 0 \end{bmatrix}, D = \begin{bmatrix} c & 0 \\ 0 & 1 \end{bmatrix},$$

where $c = \cosh \tau$, $s = \sinh \tau$, τ being a real parameter. This subgroup is contained in a (two-dimensional) subgroup A (from the Iwasawa decomposition $SU(2,2) = KAN$). The $f(z)$ below is for the (positive) square root of the conformal factor $h(z)$. The latter has been defined by our (2.1). To simplify notation, we, sometimes, use the same symbol (like z or \tilde{z} , below) to denote both an element of $E^{(2)}$ and a matrix in E_0 . The statement and the proof of the following theorem presume usage of rectangular coordinates in Euclidean E^6 : see Appendix A.

Theorem B.1. The image \tilde{z} of z in $E^{(2)}$ and the conformal factor $h(z)$ at z (under the lift of the (B.1)—transformation g) are as follows:

$$\left. \begin{aligned} \tilde{u}_{-1} + i\tilde{u}_0 &= [cu_{-1} - su + i(cu_0 - su_1)] f(z), \\ \tilde{u}_2 + i\tilde{u}_1 &= [cu_2 - su_{-1} + i(cu_1 - su_0)] f(z), \\ \tilde{u}_4 + i\tilde{u}_3 &= (u_4 + iu_3) f(z); \end{aligned} \right\} \tag{B.2}$$

$$h(z) = (f(z))^2 = \left[(cu_{-1} - su_2)^2 + (cu_0 - su_1)^2 \right]^{-1}. \quad (\text{B.3})$$

Proof. Notice that due to (A3) and (A4) from Appendix A, the formulas (B.2) correctly define the transformation on the level of E_0 (when z and \tilde{z} are matrices). To prove this first part, we use (B.1) in a straightforward way and (omitting routine details of the calculation) determine (B.2). At this stage of the proof we cannot be sure that $h(z)$ is the conformal coefficient. To prove that it is, apply the differential operator d to both sides of (B.2) in order to express

$$(d\tilde{u}_{-1})^2 + (d\tilde{u}_0)^2 + (d\tilde{u}_1)^2 - (d\tilde{u}_2)^2 - (d\tilde{u}_3)^2 - (d\tilde{u}_4)^2 \quad (\text{B.4})$$

in terms of differentials $du_{-1}, du_0, du_1, du_2, du_3, du_4$. Comparison of the obtained expression with (A5L) verifies (B.3). \square

Remark B.2. In the case considered, there is an alternative way to determine the conformal factor (B.3). It is as follows [17], Theorem 3: for a (1.3)-transformation g , the following equality holds for the conformal factor at z :

$$h(z) = \det \left(\left[A - (Az + B)(Cz + D)^{-1} C \right] z (Az + B)^{-1} \right). \quad (\text{B.5})$$

One can verify that (B.5), when applied in the (B.1)-case, results in (B.3).

It is of interest to determine all *fixed points* (that is, matrices $z = \begin{bmatrix} z_1 & z_2 \\ z_3 & z_4 \end{bmatrix}$ with $g(z) = z$ property) of the transformation (B.1).

Scholium B.3. The totality of all fixed points of (B.1) is a pair of circles. One of the circles is given by equations $z_3 = 1, z_1 = z_4 = 0$. The other circle is given by equations $z_3 = -1, z_1 = z_4 = 0$.

Proof. As it follows from (B.1), the totality of all fixed points is the solution set of

$$\begin{bmatrix} z_1 & z_2 \\ z_3 & z_4 \end{bmatrix} \begin{bmatrix} c + sz_3 & sz_4 \\ 0 & 1 \end{bmatrix} = \begin{bmatrix} z_1 & z_2 \\ cz_3 + s & z_4 \end{bmatrix}, \quad (\text{B.6})$$

equality of two matrices. Comparison of first entries in second rows results in

$$s(z_3)^2 = s. \quad (\text{B.7})$$

Since g is not an identity transformation, $(z_3)^2 = 1$. If $z_3 = 1$, then comparison of first entries in the first rows results in $(c + s)z_1 = z_1$, that is, $z_1 = 0$. Comparison of second entries in second rows results in $z_4 = 0$. Now, if $z_3 = -1$ in (B.7), then, again, $z_1 = z_4 = 0$. \square

Our next goal is to prove that each fixed point (of a given (B.1)—transformation) is an extreme point of the conformal coefficient $h(z)$: maximum is reached at each point of one circle whereas minimum is reached at each point of the other circle.

Scholium B.4. If $(u_1)^2 + (u_2)^2 < 1$ at z_0 , then z_0 is not a point of extremum for $h(z)$.

Proof. If $(u_1)^2 + (u_2)^2 < 1$ at z_0 , then u_1, u_2 can be chosen as two (of the total of four) free real variables at the vicinity of z_0 . If z_0 is a point of extremum for $h(z)$, then each of the two partial derivatives of h (w.r.t. u_1 , and w.r.t. u_2) vanish at z_0 . However, that would have resulted in vanishing of $h(z)$ at z_0 . \square

Corollary B.5. If $z = (u_{-1}, u_0, u_1, u_2, u_3, u_4)$ is an extreme point for h , then $u_3 = u_4 = 0$ (that is, $z_1 = z_4 = 0$).

Corollary B.6. At the point of extremum for the conformal coefficient, either $u_{-1}u_2 + u_0u_1 = 1$, or $u_{-1}u_2 + u_0u_1 = -1$ (that is, $z_3 = -1$, or $z_3 = 1$).

Proof follows from the expression $h(z) = \left[c^2 + s^2 - 2sc(u_{-1}u_2 + u_0u_1) \right]^{-1}$ which holds at every point z where $u_3 = u_4 = 0$. \square

Corollary B.7. An extreme value of $h(z)$ is reached at z_0 if and only if z_0 is a fixed point of (B.1). One can verify that the two extreme values are e^τ and $e^{-\tau}$ where τ is the (non-zero) value of the parameter in (B.1).

Appendix C: Evaluations of Integrals (2.8) for the Case of Appendix B Transformations in $E^{(2)}$

We start with the form

$$h = (cu_{-1} - su_2)^2 + (cu_0 - su_1)^2 \quad (C.1)$$

on the torus $T = S^1 \times S^3$, see our Theorem B.1. Now, T^+ is for the part of T where $h \geq 1$, T^- is for the part of T where $h \leq 1$. Introduce

$$I_k^+ = \int_{T^+} h^k, \quad (C.2)$$

$$I_k^- = \int_{T^-} h^k; \quad (C.3)$$

where in both cases we have in mind the volume form which has been introduced on T in Section 2.

A word of caution: the function (C.1) is the inverse of the conformal coefficient (B.3). Nevertheless, the findings (which follow) of this Appendix C are relevant to the Appendix A content since k in (C.2), (C.3) can be any integer.

The majority of these Appendix C findings are due to V. V. Ivanov (Sobolev Institute of Mathematics, Novosibirsk, Russia).

Parameterize T as follows:

$$\left. \begin{aligned} u_{-1} = \cos(\phi - \psi), u_0 = \sin(\phi - \psi), u_1 = \varrho \cos(\psi), u_2 = \varrho \sin(\psi), \\ u_3 = \sqrt{1 - \varrho^2} \cos(\xi), \quad u_4 = \sqrt{1 - \varrho^2} \sin(\xi). \end{aligned} \right\} \quad (C.4)$$

In terms of these parameters, (C.1) becomes

$$h = h(\varrho, \phi) = c^2 - 2c\varrho s \sin(\phi) + s^2 \varrho^2. \quad (C.5)$$

The integrals (C.2), (C.3) are reduced as follows:

$$I_k^\pm = 4\pi^2 J_k^\pm, \quad (C.6)$$

where

$$J_k^\pm = \iint_{\Omega^\pm} \varrho h^k(\varrho, \phi) d\varrho d\phi. \quad (C.7)$$

Here we consider the rectangle $0 \leq \varrho \leq 1, 0 \leq \phi \leq 2\pi$, and $\Omega^+ = \{(\varrho, \phi) : h \geq 1\}$, $\Omega^- = \{(\varrho, \phi) : h \leq 1\}$. Notice that the integrals (C.6) and (C.7) are independent of the sign of $s = \sinh(\tau)$ which allows us to stay with $s > 0$, only. The next step is to interpret ϱ, ϕ as polar coordinates on the x, y plane:

$$x = \varrho \cos(\phi), y = \varrho \sin(\phi). \quad (C.8)$$

Our function (C.5) becomes

$$h(x, y) = (sx)^2 + (c - sy)^2, \quad (C.9)$$

whereas Ω^+, Ω^- are to be converted into D^+, D^- with their union being the unit disc D centered at the origin $(0,0)$ of the x, y plane. Finally, introduce coordinates r, α :

$$x = r \sin(\alpha), y = c/s - r \cos(\alpha). \quad (C.10)$$

r being the distance between $P = (0, c/s)$ and $Q = (x, y)$, whereas the angle α , in radians, is an angle between vectors $\{0, -1\}$ and PQ . Expression (C.9) becomes

$$h(r, \alpha) = s^2 r^2. \quad (C.11)$$

Introduce an (acute) angle ω which is determined by any of the relations

$$\omega \cos(\omega) = 1/c, \quad \sin(\omega) = s/c, \quad \tan(\omega) = s. \quad (C.12)$$

Omitting a few more (straightforward) technicalities, we obtain

$$J_k^\pm = (\pm 2/s^2) \int_0^\omega \int_1^{H_\pm(\alpha)} r^{2k+1} dr d\alpha. \quad (\text{C.13,C14})$$

The upper limits $H_+(\alpha)$, $H_-(\alpha)$ are as follows:

$$H_\pm(\alpha) = c \cos(\alpha) \pm \sqrt{c^2 \cos^2 \alpha - 1}. \quad (\text{C.15,C.16})$$

Let us conclude in terms of the following statements.

Theorem C.1. For k not equal -1 , the integrals (C.2), (C.3) can be evaluated as follows:

$$I_k^\pm = \frac{\pm 4\pi^2}{(k+1)s^2} \int_0^\omega (H^{\pm 2k+2} - 1) d\alpha. \quad (\text{C.17})$$

For $k = -1$

$$I_{-1}^\pm = \frac{\pm 8\pi^2}{s^2} \int_0^\omega \ln(H_\pm(\alpha)) d\alpha. \quad (\text{C.18})$$

Theorem C.2. For every integer k ,

$$I_{-k-2}^+ = I_k^-, \quad I_{-k-2}^- = I_k^+. \quad (\text{C.19})$$

Theorem C.3. For a nonnegative k , each of the integrals (C.17) is a finite linear combination of integrals A_m , B_m where

$$\left. \begin{aligned} A_m &= \int_0^\omega c^{2m} \cos^{2m}(\alpha) d\alpha, \\ B_m &= \int_0^\omega c^{2m-1} \cos^{2m-1}(\alpha) \sqrt{c^2 \cos^2(\alpha) - 1} d\alpha. \end{aligned} \right\} \quad (\text{C.20})$$

Remark C.4. Each of the integrals (C.20) is an elementary one and it can be expressed as a polynomial in s and ω .

Recall notations a, b, c, d of Section 2 (see the line prior to Formula (2.9)) for the integrals which are of our utmost interest.

Theorem C.5. The integrals a, b, c, d , are as follows:

$$\begin{aligned} a &= \left\{ 4\pi^2 / [3(s^2)] \right\} \{ 32A_3 - 48A_2 + 18A_1 - 2A_0 + 32B_3 - 32B_2 + 6B_1 \}, \\ b &= 4\pi^2 \left[\pi/2 + \omega + (s - \omega) / (s^2) \right], \\ c &= \left\{ 4\pi^2 / [3(s^2)] \right\} \{ 32A_3 - 48A_2 + 18A_1 - 2A_0 - 32B_3 + 32B_2 - 6B_1 \}, \\ d &= 4\pi^2 \left[\pi/2 - \omega - (s - \omega) / (s^2) \right]. \end{aligned}$$