RELATIVITY WITHOUT TEARS

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Special relativity is no longer a new revolutionary theory but a firmly established cornerstone of modern physics. The teaching of special relativity, however, still follows its presentation as it unfolded historically, trying to convince the audience of this teaching that Newtonian physics is natural but incorrect and special relativity is its paradoxical but correct amendment. I argue in this article in favor of logical instead of historical trend in teaching of relativity and that special relativity is neither paradoxical nor correct (in the absolute sense of the nineteenth century) but the most natural and expected description of the real space-time around us valid for all practical purposes. This last circumstance constitutes a profound mystery of modern physics better known as the cosmological constant problem.

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Preface

“To the few who love me and whom I love — to those who feel rather than to those who think — to the dreamers and those who put faith in dreams as in the only realities — I offer this Book of Truths, not in its character of Truth-Teller, but for the Beauty that abounds in its Truth; constituting it true. To these I present the composition as an Art-Product alone; let us say as a Romance; or, if I be not urging too lofty a claim, as a Poem.

What I here propound is true: — therefore it cannot die: — or if by any means it be now trodden down so that it die, it will rise again ‘to the Life Everlasting’.

Nevertheless, it is a poem only that I wish this work to be judged after I am dead” [1].

(811)
1. Introduction

The emergence of special relativity constitutes one of the major revolutions in physics. Since then this revolution is repeated over and over in physics students' minds. “Often the result is to destroy completely the confidence of the student in perfectly sound and useful concepts already acquired” [2]. What is the problem the students stumble upon?

Maybe the following passage from [3] might give you an insight: “At first, relativity was considered shocking, anti-establishment and highly mysterious, and all presentations intended for the population at large were meant to emphasize these shocking and mysterious aspects, which is hardly conducive to easy teaching and good understanding. They tended to emphasize the revolutionary aspects of the theory whereas, surely, it would be good teaching to emphasize the continuity with earlier thought”.

The standard presentation of special relativity follows the royal way Einstein laid down in his breakthrough paper [4] and is based on two postulates. The first postulate is the Principle of Relativity that the laws of physics are the same in all inertial reference frames. This is precisely what the students are taught in their undergraduate mechanics course and should hardly create any conceptual problems because the Principle of Relativity “appeals to our common sense with irresistible force” [5].

There is a subtle difference in the formulations of the second postulate as given by Einstein and as presented in modern textbooks [6]. In Einstein’s paper [4] the postulate states that “any ray of light moves in the ‘stationary’ system of coordinates with the determined velocity $c$, whether the ray be emitted by a stationary or by a moving body”. Taken separately, this statement is also both orthodox and obvious in the context of the pre-1905 physics with luminiferous æther as its conceptual basis. Indeed, after a light wave is launched in the æther it propagates with the speed which is completely independent of the state of motion of the light source, determined solely by the elastic properties of the æther.

It is the combination of these two postulates, considered impossible in the pre-1905 physics, that shattered the very foundations of contemporary physics and changed forever our perspective of space-time.

Many modern textbooks make explicit what is so cunningly hidden in Einstein’s seemingly harmless postulates. They state the second postulate as follows: “(in vacuum) light travels rectilinearly at speed $c$ in every direction in every inertial frame” [7, 8]. And here is buried the root of confusion students experience while studying special relativity. Because for Newtonian intuition “to take as a postulate that the speed of light is constant relative to changes in reference frame is to assume an apparent absurdity. It goes against common sense. No wonder, thinks a student, that we can derive
other absurdities, such as time dilation and length contraction, from the premises” [6].

Laying aside the question of culture shock students experience when confronting the second postulate, there are several reasons, from the modern perspective, not to base relativity on the second postulate.

First of all, it anchors special relativity in the realm of electromagnetism while now we know that relativity embraces all natural phenomena. Strong and weak interactions, unknown to Einstein in 1905, are also Lorentz invariant.

The second postulate assumes light propagation in vacuum. But which vacuum? Vacuum in modern physics is quite different from just empty space and looks, in a sense, like an incarnation of æther, which “at present, renamed and thinly disguised, dominates the accepted laws of physics” [9]. Of course this new “æther”, being Lorentz invariant and quantum mechanical, has nothing in common with the æther of the nineteenth century. Anyway, in some cases it reveals itself as an nontrivial medium capable to alter the propagation properties of light.

For example, between conducting plates [10,11] or in a background gravitational field [12,13] light can propagate with speeds greater than c. Yet the Lorentz invariance remains intact at a fundamental level [14]. Simply the boundary conditions or the background field single out a preferred rest frame and the ground state (quantum vacuum) becomes not Lorentz-invariant. The presence of such not Lorentz-invariant “æther” can be detected in contrast to the situation in infinite space with no boundaries [14].

Whether light propagates with invariant velocity c is subject of photon being massless. This masslessness of the photon by itself originates from the particular pattern of the electroweak symmetry breaking. However, there is no compelling theoretical reason for the photon to be strictly massless: a tiny photon mass would not destroy the renormalizability of quantum electrodynamics and hence the beautiful agreement between its predictions and experiment [15]. Moreover, it was shown that the photon mass can be generated by inflation [16]. As the current universe is entering another phase of inflation, according to the supernovae results [17], the photon should have a miniscule mass [16] of about $10^{-42}$ GeV$/c^2$, far below of the present experimental limits. Anyway this miniscule mass makes the photon not the best choice the special relativity to base on.

Surprisingly, it was known for a long time that in fact one does not need the second postulate to develop the special relativity. The Relativity Principle alone along with some “self-evident” premises such as the homogeneity of space and time and the isotropy of space would suffice.
To my knowledge, von Ignatowsky [18, 19] was the first to discover this remarkable fact as early as in 1910, followed by Frank and Rothe [20]. “However, like numerous others that followed these papers have gone largely unnoticed” [21]. An impressive list of those largely neglected and forgotten papers can be compiled from citations in [14, 22, 23].

The idea has got a better chance after it was rediscovered in the modern context [24–26]. At least, it attracts interest up to date [27–31]. “At this point we can sharpen and somewhat displace some typical questions of historical character: Why, although it was logically and epistemologically perfectly possible, did not Einstein see the connection between his two postulates? Was his embarrassment a kind of a subconscious inkling of this connection? Why such early papers as that of Ignatowsky did not catch more the interest of physicists, historians and philosophers of science? Why such a connection is rarely mentioned in pedagogical presentations of Relativity?” [32].

I think answers to these questions can be found partly in the concrete historical circumstances of the emergence and development of relativity and partly in the “intellectual inertia” [33] of society.

As was already mentioned, in Einstein’s original form the second postulate was not shocking at all for contemporary physicists. The focus was displaced to the weird but logically inevitable consequences that followed when the second postulate was combined with the Relativity Principle. The one-postulate derivations of Ignatowsky et al. involved a somewhat higher level of mathematics (group theory, more intricate and rather less familiar analysis; See for example [25, 26], or earlier considerations by Lalan [34]). I suspect, this was considered as an unnecessary complication, making the approach “unavailable for a general education physics course” [26].

The focus has shifted since then from reference frames and clock synchronization to symmetries and space-time structure, and the situation is different today. For contemporary students the luminiferous æther is just an historical anachronism and can not serve as the epistemological basis for the second postulate. Einstein’s brilliant magic when he, “having taken from the idea of light waves in the æther the one aspect that he needed” [6], declared later in his paper that “the introduction of a ‘luminiferous æther’ will prove to be superfluous” [4], does not work any more. Therefore Ignatowsky’s approach is much more appealing today than it was in 1910, because it leads to Lorentz transformations, which are at the heart of special relativity, “without ever having to face the distracting sense of paradox that bedevils more conventional attempts from the very first steps” [26].

Below I will try to show that, combining ideas from [24–26, 35], it is possible to make the one-postulate derivation of Lorentz transformations mathematically as simple as was Einstein’s original presentation.
In fact, much richer algebraic and geometric structures are lurking behind special relativity [36]. Some of them will be considered in subsequent chapters.

2. Relativity without light

Let an inertial frame of reference $S'$ move along the $x$-axis with velocity $V$ relative to the "stationary" frame $S$.

A simple glance at Fig. 1 is sufficient to write down the Galilean transformation that relates the $x$-coordinates of some event (for example, an explosion) in the frames $S$ and $S'$,

$$x = V t + x'.$$

But in (1) we have implicitly assumed that meter sticks do not change their lengths when gently set in uniform motion. Although intuitively appealing, this is not quite obvious. For example, according to Maxwell's equations, charges at rest interact only through the Coulomb field while in motion they experience also the magnetic interaction. Besides, when a source moves very fast its electric field is no longer spherically symmetrical. It is therefore not unreasonable to expect that a meter stick set in rapid motion will change shape in so far as electromagnetic forces are important in ensuring the internal equilibrium of matter [2]. Anyway we just admit this more general possibility and change (1) to

$$x = V t + k(V^2) x',$$

where the scale factor $k(V^2)$ accounts for the possible change in the length of the meter stick.
The Relativity Principle and the isotropy of space are implicit in (2) because we assumed that the scale factor depends only on the magnitude of the relative velocity $V$.

Equation (2) allows us to express the primed coordinates through the unprimed ones

$$x' = \frac{1}{k(V^2)} (x - Vt).$$

(3)

Then the Relativity Principle tells us that the same relation holds if unprimed coordinates are expressed through the primed ones, with $V$ replaced by $-V$, because the velocity of $S$ with respect to $S'$ is $-V$. Therefore,

$$x = \frac{1}{k(V^2)} (x' + Vt') = \frac{1}{k(V^2)} \left[ \frac{1}{k(V^2)} (x - Vt) + Vt' \right].$$

Solving for $t'$, we get

$$t' = \frac{1}{k(V^2)} \left[ t - \frac{1 - k^2(V^2)}{V} x \right].$$

(4)

We see immediately that time is not absolute if the scale factor $k(V^2)$ is not identically one.

From (3) and (4) we can infer the velocity addition rule

$$v'_x = \frac{dx'}{dt'} = \frac{dx - Vdt}{dt - \frac{1 - k^2}{V} dx} = \frac{v_x - V}{1 - \frac{1 - k^2}{V} v_x}.$$

In what follows it will be convenient to write down the velocity addition rule with unprimed velocity expressed through the primed one,

$$v_x = \frac{v'_x + V}{1 + \frac{1 - k^2}{V} v'_x} \equiv F(v'_x, V).$$

(5)

If we change the signs of both velocities $v'_x$ and $V$ it is obvious that the sign of the resulting velocity $v_x$ will also change. Therefore $F$ must be an odd function of its arguments

$$F(-x, -y) = -F(x, y).$$

(6)

Consider now three bodies $A$, $B$ and $C$ in relative motion. Let $V_{AB}$ denote the velocity of $A$ with respect to $B$. Then

$$V_{BA} = -V_{AB}.$$

(7)

This is the reciprocity principle already used above. I think it is obvious enough to require no proof in an introductory course of special relativity.
However, in fact it can be deduced from the Relativity Principle, space-time homogeneity and space isotropy [23, 37]. Some subtleties of such proof are discussed in [38].

Now, using (6) and (7), we get [26]

\[ F(V_{CB}, V_{BA}) = V_{CA} = -V_{AC} \]

\[ = -F(V_{AB}, V_{BC}) = -F(-V_{BA}, -V_{CB}) = F(V_{BA}, V_{CB}). \]

Therefore \( F \) is a symmetric function of its arguments. Then \( F(v'_x, V) = F(V, v'_x) \) immediately yields, according to (5),

\[ \frac{1 - k^2(V^2)}{V} v'_x = \frac{1 - k^2(v'_x^2)}{v'_x^2} V, \]

or

\[ \frac{1 - k^2(V^2)}{V^2} = \frac{1 - k^2(v'_x^2)}{v'_x^2} \equiv K, \quad (8) \]

where at the last step we made explicit that the only way to satisfy equation (8) for all values of \( V \) and \( v'_x \) is to assume that its left- and right-hand sides are equal to some constant \( K \). Then the velocity addition rule will take the form

\[ v_x = \frac{v'_x + V}{1 + K v'_x V}. \quad (9) \]

If \( K = 0 \), one recovers the Galilean transformations and velocity addition rule \( v_x = v'_x + V \).

If \( K < 0 \), one can take \( K = -\frac{1}{c^2} \) and introduce a dimensionless parameter \( \beta = \frac{V}{c} \). Then

\[ x' = \frac{1}{\sqrt{1 + \beta^2}} (x - Vt), \]

\[ t' = \frac{1}{\sqrt{1 + \beta^2}} \left( t + \frac{V}{c} x \right), \quad (10) \]

while the velocity addition rule takes the form

\[ v_x = \frac{v'_x + V}{1 - \frac{v'_x V}{c^2}}. \quad (11) \]

If \( v'_x = V = \frac{c}{2} \), then (11) gives \( v_x = \frac{4}{3} c \). Therefore velocities greater than \( c \) are easily obtained in this case. But if \( v'_x = V = \frac{4}{3} c \), then

\[ v_x = \frac{\frac{8}{3} c}{1 - \frac{16}{9}} = \frac{-24}{7} c < 0. \]
Therefore two positive velocities can sum up into a negative velocity! This is not the only oddity of the case $K < 0$. For example, if $v'_x V = c^2$, then $v'_x$ and $V$ will sum up into an infinite velocity.

If we perform two successive transformations $(x, t) \to (x', t') \to (x'', t'')$, according to (10) with $\beta = 4/3$, we end up with

$$x'' = -\frac{7}{25} x - \frac{24}{25} c t.$$ 

This can not be expressed as the result of a transformation $(x, t) \to (x'', t'')$ of type (10) because the coefficient $1/\sqrt{1 + \beta^2}$ of $x$ in (10) is always positive. Hence the breakdown of the group law.

However, the real reason why we should discard the case $K < 0$ is the absence of causal structure. The transformations (10) can be recast in the form

$$x'_0 = \cos \theta x_0 + \sin \theta x,$$
$$x' = \cos \theta x - \sin \theta x_0,$$

(12)

where $x_0 = ct$ and $\cos \theta = \frac{1}{\sqrt{1 + \beta^2}}$. This is the usual rotation and hence the only invariant quantity related to the event is $x_0^2 + x^2$. This Euclidean norm does not allow us to define an invariant time order between events, much like of the ordinary three-dimensional space where one can not say which points precede a given one.

Finally, if $K = \frac{1}{c^2} > 0$, one gets the Lorentz transformations

$$x' = \frac{1}{\sqrt{1 - \beta^2}} (x - Vt),$$
$$t' = \frac{1}{\sqrt{1 - \beta^2}} \left( t - \frac{V}{c^2} x \right),$$

(13)

with the velocity addition rule

$$v_x = \frac{v'_x + V}{1 + \frac{v'_x V}{c^2}}.$$ 

(14)

Now $c$ is an invariant velocity, as (14) shows. However, in the above derivation nothing as yet indicates that $c$ is the velocity of light. Only when we invoke the Maxwell equations it becomes clear that the velocity of electromagnetic waves implied by these equations must coincide with $c$ if we want the Maxwell equations to be invariant under Lorentz transformations (13).
The choice between the Galilean ($K = 0$) and Lorentzian ($K > 0$) cases is not only an experimental choice. Some arguments can be given supporting the idea that this choice can be made even on logical grounds.

First of all, special relativity is more friendly to determinism than the classical Newtonian mechanics [39]. An example of indeterministic behavior of a seemingly benign Newtonian system is given by Xia’s five-body supertask [40, 41].

Two symmetrical highly eccentric gravitationally bound binaries are placed at a large distance from each other, while a fifth body of much smaller mass oscillates between the planes of these binaries (see Fig. 2). It can be shown [40] that there exists a set of initial conditions under which the binaries from this construction will escape to spatial infinity in a finite time, while the fifth body will oscillate back and forth with ever increasing velocity.

The time reverse of Xia’s supertask is an example of “space invaders” [39] — particles appearing from spatial infinity without any causal reason. This indeterminacy of idealized Newtonian world is immediately killed by special relativity because without unbounded velocities there are no space invaders.

That causality and Lorentz symmetry are tightly bound is a remarkable fact. Another less known royal way to relativity can be traced back to Alfred Robb’s synthetic approach [42] which emphasizes the primary role of causal order relation in determining the space-time geometry. Later this line of reasoning was further developed by Reichenbach [43] from the side of philosophy and Alexandrov [44] and Zeeman [45] from the side of mathematics. The famous Alexandrov–Zeeman theorem states that, in a sense, causality implies Lorentz symmetry.
Group theory provides another argument to prefer the Lorentzian world over the Galilean one \cite{46} because mathematically the Lorentz group $G_c$ has much simpler and elegant structure than its singular limit $G_\infty$ which is the symmetry group of the Galilean world. This argument goes back to Minkowski. In his own words, cited in \cite{46}, from his famous Cologne lecture Raum und Zeit:

“Since $G_c$ is mathematically more intelligible than $G_\infty$, it looks as though the thought might have struck some mathematician, fancy free, that after all, as a matter of fact, natural phenomena do not possess an invariance with the group $G_\infty$, but rather with the group $G_c$, $c$ being finite and determinate, but in ordinary units of measure extremely great. Such a premonition would have been an extraordinary triumph for pure mathematics. Well, mathematics, though it can now display only staircase-wit, has the satisfaction of being wise after the event, and is able, thanks to its happy antecedents, with its senses sharpened by an unhampered outlook to far horizons, to grasp forthwith the far-reaching consequences of such a metamorphosis of our concept of nature”.

3. Relativistic energy and momentum

Although Minkowski’s contribution and his notion of space-time is the most crucial event placing relativity in the modern context, students reluctantly swallow the four-vector formalism due to their Newtonian background. Therefore, some elementary derivation of the relativistic expressions for energy and momentum, stressing not the radical break but continuity with concepts already acquired by students \cite{2}, is desirable in a general education physics course.

Usually that is done by using the concept of relativistic mass — another unfortunate historical heritage stemming from an inappropriate generalization of the Newtonian relationship between momentum and mass $\vec{p} = m\vec{v}$ to the relativistic domain.

Crystal-clear arguments were given \cite{47,48} against the use of relativistic mass as a concept foreign to the logic of special relativity. However, the “intellectual inertia” still prevails, the most famous (wrong!) formula associated to special relativity for the general public is still $E = mc^2$, and students are still exposed to the outmoded notion of velocity dependent mass (unfortunately, this is maybe the only concept of special relativity they absorb with ease because of its Newtonian roots).

In fact, even to ensure continuity with Newtonian concepts there is no need to preserve historical artifacts in teaching special relativity a hundred years after its creation. Below a modification of Davidon’s derivation \cite{22,49} of relativistic energy and momentum is presented which quite gently harmonizes with students’ Newtonian background.
Suppose a ball cools down by emitting photons isotropically in its rest frame $S'$. Let the total energy emitted during some time be $E'$. Because of isotropy, the total momentum carried away by radiation is zero and the ball will stay motionless. The Principle of Relativity then implies that in every other inertial frame the velocity of the ball is also unchanged.

Let us see how things look in the reference frame $S$ in which the ball moves with velocity $V$ along the $x$-axis. A bunch of photons emitted within the solid angle $d\Omega' = 2\pi \sin \theta' d\theta'$ around the polar angle $\theta'$ in the frame $S'$ has total energy $dE' = E' d\Omega'$ in this frame. Due to the Doppler effect, the corresponding energy in the frame $S$ will be

$$dE = \gamma dE' (1 + \beta \cos \theta') = -\gamma E' \frac{1}{2} (1 + \beta \cos \theta') d \cos \theta'.$$

(Note that the Doppler formula for the photon frequency shift follows from the Lorentz transformations quite easily [7]. We have also assumed the relation $E = \hbar \omega$ for the photon energy from elementary quantum theory.)

Therefore, the total energy emitted in the frame $S$ is

$$E = -\gamma E' \frac{1}{2} \int_0^\pi (1 + \beta \cos \theta') d \cos \theta' = \gamma E'. $$

In the frame $S$, the radiation is no longer isotropic and therefore it takes away some momentum. Let us calculate how much. The emission angle of the bunch of light in the frame $S$ can be found from the aberration formula (which follows from the velocity addition rule (14))

$$\cos \theta = \frac{\cos \theta' + \beta}{1 + \beta \cos \theta'}. $$

But each photon of energy $\hbar \omega$ carries the momentum $\hbar \omega / c$ (another elementary fact from the quantum theory). Therefore the $x$-component of the momentum of the light-bunch is

$$dP_x = \frac{dE}{c} \cos \theta = -\frac{\gamma E'}{2c} (\cos \theta' + \beta) d \cos \theta'. $$

Integrating, we get the total momentum taken away by radiation in the $x$-direction

$$P_x = -\frac{\gamma E'}{2c} \int_0^\pi (\cos \theta' + \beta) d \cos \theta' = \gamma \beta \frac{E'}{c}. $$

Therefore, for the energy and momentum of the ball in the frame $S$ we should have

$$\Delta E = -\gamma E', \quad \Delta p = -\gamma \beta \frac{E'}{c}. \quad (15)$$
It is natural to assume that the momentum and velocity of the ball are parallel to each other,
\[ \vec{p} = N(V)\vec{V}, \quad (16) \]
where \( N(V) \) is some unknown function (let us call it the inertia of the ball \([22, 49]\)). Then
\[ \Delta p = V\Delta N, \]
because \( \Delta V = 0 \) as explained above. But (15) implies that \( \Delta p = \frac{\beta \Delta E}{c} \). Therefore, \( \frac{\beta \Delta E}{c} = V\Delta N \) and
\[ \Delta E = c^2\Delta N. \quad (17) \]

Although this result was obtained in the particular circumstances, we now assume that it is universally valid. That is, we assume that every change in the energy of a body implies the corresponding change in its inertia according to (17). If the body is subject to a force \( \vec{F} \) then
\[ \frac{dE}{dt} = \vec{F} \cdot \vec{V} = \vec{V} \cdot \frac{d\vec{p}}{dt} = \vec{V} \cdot \left( V^2\frac{dN}{dt} + N \frac{dV}{dt} \right) = V^2\frac{dN}{dt} + \frac{N}{2} \frac{dV^2}{dt}. \]
Using (17) we get
\[ c^2\frac{dN}{dt} = V^2\frac{dN}{dt} + \frac{N}{2} \frac{dV^2}{dt}, \]
or
\[ \frac{dN}{N} = \frac{dV^2}{2(c^2 - V^2)}. \]

Integrating, we get
\[ N = \frac{N_0}{\sqrt{1 - \beta^2}}, \quad (18) \]
where \( N_0 \) is an integration constant. Therefore,
\[ \vec{p} = \frac{N_0\vec{V}}{\sqrt{1 - \beta^2}} = \frac{m\vec{V}}{\sqrt{1 - \beta^2}} \quad (19) \]
because by considering the nonrelativistic limit \( \beta \ll 1 \) we conclude that \( N_0 \) is just the mass of the body. It follows from (17) and (18) that
\[ E = Nc^2 = \frac{mc^2}{\sqrt{1 - \beta^2}} \quad (20) \]
up to irrelevant overall constant. This derivation of relativistic energy and momentum clarifies the real meaning of the notorious $E = mc^2$. It is the inertia of the body, defined as the coefficient of the velocity in the expression of the momentum $\vec{p} = N\vec{V}$, to which the energy of the body is proportional, while the mass of the body is an invariant, frame-independent quantity — much like the Newtonian concept of mass.

4. Relativity without reference frames

You say that relativity without reference frames does not make even a linguistic sense? Yes, of course, if the words are understood literally. But “the name ‘theory of relativity’ is an unfortunate choice. Its essence is not the relativity of space and time but rather the independence of the laws of nature from the point of view of the observer” [50]. Of course, I do not advocate changing the name, so celebrated, of the theory. But it should be instructive to bear in mind the conventionality of the name and the shift of focus from the length contraction and time dilation (which are rather obvious effects in Minkowski’s four-dimensional geometric formulation) to space-time structure and symmetries.

Above we have discarded the $K < 0$ possibility as unphysical. But should we really? The main argument was that the corresponding Euclidean space-time does not support causal structure and, therefore, it is in fact some kind of a timeless nirvana. The real turbulent universe around us is certainly not of this type. But we cannot exclude that it may contain inclusions of Euclidean domains, maybe formed at the centers of black holes as a result of quantum signature change during gravitational collapse [51]. Moreover, the Hartle–Hawking’s ‘No Boundary Proposal’ in quantum cosmology assumes that Euclidean configurations play an important role in the initial wavefunction of the Universe [52]. Loosely speaking, it is suggested that the whole Universe was initially Euclidean and then by quantum tunneling a transition to the usual Lorentzian space-time occurred.

But once we accept the Euclidean space-time as physically viable we have to change the principal accents of the formalism, depriving reference frames of their central role, because in timeless Euclidean nirvana there are no observers and hence no reference frames. The Principle of Relativity then should be replaced by the more general concept of symmetry transformations which leave the space-time structure invariant. In doing so, a question naturally arises as to what is the concept of space-time geometry which embraces even space-times inhospitable to intelligent life. And at this point it is a good thing to acquaint students with the Erlangen program of Felix Klein.
5. What is geometry?

The word *geometry* is derived from Greek and means *earth measurement*. It is not surprising, therefore, that beginning from Gauss and Riemann the length concept is considered to be central in geometry. The distance between two infinitesimally close points in Riemannian geometry is given by some positive definite, or at least non-degenerate, quadratic differential form

\[ ds^2 = \sum_{i,j=1}^{n} g_{ij}(x) \, dx^i dx^j. \]

Other geometric concepts like the angle between two intersecting curves can be defined in terms of the metric \( g_{ij} \).

The main feature of the Riemannian geometry is that it is local and well suited for field theory in physics [53], general relativity being the most notable example. However, some interesting and important questions are left outside the scope of Riemannian geometry.

For example, there are two basic geometric measurements: the determination of the distance between two points and the determination of the angle between two intersecting lines. The corresponding measures are substantially different [54]. The distance between two points is an algebraic function of their coordinates and therefore it is uniquely defined (up to a sign). In contrast, the angle between two intersecting lines is determined as a transcendental (trigonometric) function of coordinates and is defined only up to \( 2\pi n \), for an arbitrary integer \( n \). As a result, every interval can easily be divided into an arbitrary number of equal parts, while the division of an arbitrary angle into three equal angles by using only a compass and a straightedge is an ancient problem proved impossible by Wantzel in 1836.

This difference between length and angle measurements is just implied in Riemannian geometry without any explanation and hardly disturbs our intuition until we recognize its Euclidean roots. The existence of non-Euclidean geometries when calls for a more careful examination of the question how the corresponding measures arise. For this purpose we should take another road to geometry, the Kleinian road where the equality of figures is a basic geometrical concept [55].

In Euclidean geometry the equality of two figures means that they can be superimposed with an appropriate rigid motion. The axiom that if \( A \) equals \( B \) and \( B \) equals \( C \), then \( A \) equals \( C \), in fact indicates [55] that rigid motions form a group. A far-reaching generalization of this almost trivial observation was given by Felix Klein in his famous inaugural lecture [56] that he prepared as professor at Erlangen in 1872 but never actually gave. In this lecture Klein outlined his view of geometry which later became known as the Erlangen program. Actually the Erlangen program is just the basic
principle of Galois theory applied to geometry [57]. According this principle discovered by Galois around 1832 one can classify the ways a little thing (a point or figure) can sit in a bigger thing (space) by keeping track of the symmetries of the bigger thing that leave the little thing unchanged [57].

Every particular geometry determines the group of transformations (motions) which preserve this geometry. Klein calls this group the principal group and for him the geometry is just the study of invariants of the principal group because genuine geometric properties of a figure are only those which remain unchanged under transformations belonging to the principal group.

The converse is also true: a group $G$ acting on a space $X$ determines some geometry [58]. In classical geometries all points of $X$ look alike (the space is homogeneous). In the group-theoretical language this means that $G$ acts on $X$ transitively; that is, for every pair of points $x$ and $y$ of $X$ there exists a symmetry transformation (motion) $g \in G$ which takes $x$ into $y$: $y = g(x)$. Let $H$ be the stabilizer of a point $x$ — the set of all transformations in $G$ which leave $x$ invariant. Then $X$ can be identified in fact with the coset space $G/H$. Indeed, we have a one-to-one map of $X$ onto $G/H$ which takes each point $g(x) \in X$ to the equivalence class $[g] \in G/H$. The symmetries $x \mapsto s(x), s \in G$ of the space $X$ are then represented by the transformations $[g] \rightarrow [sg]$ of the coset space $G/H$.

Hence, interestingly, the Erlangen program provides a very Kantian view of what space (space-time) is. Kant in his *Critique of Pure Reason* denies the objective reality of space and time, which for him are only forms in which objects appear to us due to the hardwired features of our consciousness (intuition) and not the properties of objects themselves.

The Kantian character of relativity theory was advocated by Kurt Gödel [59, 60]. In Gödel’s view the relativity of simultaneity deprives time of its objective meaning and “In short, it seems that one obtains an unequivocal proof for the view of those philosophers who, like Parmenides, Kant, and the modern idealists, deny the objectivity of change and consider change as an illusion or an appearance due to our special mode of perception” [60].

One can argue that the Minkowski space-time represents a kind of reality that comes for the change to the Newtonian absolute time. In Minkowski’s own words: “Henceforth space by itself, and time by itself are doomed to fade away into mere shadows, and only a kind of union of the two will preserve an independent reality” [61].

However, the Erlangen program, if pushed to its logical extreme, indicates that the Minkowski space-time is just a useful parameterization of the coset space $G/H$. The symmetry group $G$ (the Poincaré group) and its subgroup $H$ (the Lorentz group) as the stabilizer of “points” (space-time events) are all that really does matter.
However, a too abstract group-theoretical approach is not always instructive because even if space and time are really illusions they proved to be very useful concepts “due to our special mode of perception”. Therefore we leave the interesting question of the reality of space-time to philosophers; below we follow Klein’s more classical presentation [54] because it emphasizes not the radical break but continuity with concepts already acquired by students [2].

6. Projective metrics

Our special mode of perception, especially vision, determines our implicit belief that everything is made of points as the most basic structural elements. For a blind man, however, who examines things by touching them, the most basic geometric elements are, perhaps, planes [62]. This duality between various basic geometrical elements is most naturally incorporated into projective geometry [63, 64], which makes it a good starting point for studying of different kinds of linear and angular measures [54, 65]. For simplicity we will consider only plane projective geometry to demonstrate basic principles, a generalization to higher dimensional spaces being rather straightforward.

Imagine a Euclidean plane $\mathbb{R}^2$ embedded into the three dimensional Euclidean space $\mathbb{R}^3$ equipped with a Cartesian coordinate system. The coordinate system can be chosen in such a way that the equation of the plane becomes $z = a \neq 0$. Then every point in the plane $\mathbb{R}^2$ together with the origin $(0, 0, 0)$ of the coordinate system determines a line in $\mathbb{R}^3$. Therefore points in the plane $\mathbb{R}^2$ can be considered as the remnants in this plane of the corresponding lines and can be uniquely determined by the coordinates $(x, y, z)$ of any point except $(0, 0, 0)$. Two sets $(x_1, y_1, z_1)$ and $(x_2, y_2, z_2)$ of coordinates represent the same point of $\mathbb{R}^2$ if and only if $\frac{x_1}{x_2} = \frac{y_1}{y_2} = \frac{z_1}{z_2}$. Therefore the points of $\mathbb{R}^2$ are in fact the equivalence classes of triples $[x] = (\lambda x_1, \lambda x_2, \lambda x_3)$ of real numbers with $\lambda \neq 0$. If we now define the projective plane $P^2$ as the set of all such equivalence classes except $[0] = (0, 0, 0)$, then $P^2$ will not be just the Euclidean plane $\mathbb{R}^2$ because it will contain points $(\lambda x_1, \lambda x_2, 0)$ corresponding to the lines in $\mathbb{R}^3$ which are parallel to the plane $\mathbb{R}^2$ in the Euclidean sense. These points are, of course, points at infinity from the Euclidean perspective, lying on the circle of the infinite radius. Alternatively we can consider this circle as the line at infinity in $P^2$. Other lines in $P^2$ can be identified with remnants of the planes in $\mathbb{R}^3$ incident to the origin. Such planes are uniquely determined by their normal vectors $(u_1, u_2, u_3)$, or more precisely by equivalence classes $[u] = (\lambda u_1, \lambda u_2, \lambda u_3)$, $\lambda \neq 0$. Therefore, we can consider $(u_1, u_2, u_3)$ as the (projective) coordinates of a line in $P^2$. 
Now we can forget the Euclidean scaffolding of our model of the projective plane $P^2$. Then all points as well as all lines in $P^2$ will look alike. Every pair of points $x$ and $y$ are incident to a unique line $x \times y$, and every pair of lines $\xi$ and $\eta$ are incident to the unique point $\xi \times \eta$ of their intersection. Here we assume that points and lines are given by their three projective coordinates and $\times$ denotes the usual vector product of 3-dimensional vectors.

It is important to note that projective geometry, and the projective coordinates of points and lines in it, can be defined synthetically by a small set of axioms without any reference to the Cartesian coordinates and the Euclidean concept of length [66, 67]. At that, every four points $P_1, P_2, P_3$ and $E$ such that no three of them are collinear uniquely determine a projective coordinate system [64]. In this coordinate system the basic points $P_1, P_2, P_3, E$ have coordinates

$$P_1 = (1, 0, 0), \quad P_2 = (0, 1, 0), \quad P_3 = (0, 0, 1) \text{ and } E = (1, 1, 1).$$

The projective coordinates of a given point in two different coordinate systems are related by a projective transformation

$$x'_i = \sum_{i=1}^{3} A_{ij}x_j, \quad \text{det} (A_{ij}) \neq 0. \quad (21)$$

Projective transformations form a group and projective geometry studies invariants of this group. In particular, incidence relations are invariant under projective transformations and constitute the most basic geometric notion in projective geometry.

A point $x$ is incident to a line $\xi$ if and only if the scalar product $x \cdot \xi$ is equal to zero. Remarkable symmetry of this condition reflects the duality property of the projective plane: the notions ‘line’ and ‘point’ can be used interchangeably in the plane projective geometry. Every figure in $P^2$ can be considered on equal footing as made from points or as made from lines. This duality between points and lines is surprising. Our intuition does not grasp it, as our geometric terminology witnesses [64].

Another important projective invariant is the cross-ratio of four collinear points. If points $x, y, z, z'$ are collinear, so that $z = \lambda x + \mu y$ and $z' = \lambda' x + \mu' y$, then the cross-ratio of these points equals

$$R(x, y, z, z') = \frac{\mu\lambda'}{\lambda\mu'}.$$

The notions ‘angle’ and ‘distance’ are not projective geometry notions because angles and distances are not invariant under projective transformations. Therefore, to define these notions we should select a subgroup of the
projective group that will leave invariant appropriately defined angles and distances. Selecting different subgroups, we get different geometries. In the spirit of the Erlangen program all these geometries are equally feasible. It remains to find a “natural” definition of angles and distances for a given subgroup, not wildly different from what we intuitively expect from these kinds of measurements.

Both measurements, be it the measurement of an angle or distance, share common features [54]. An unknown distance is compared to some standard length regarded as the unit length. This comparison in fact involves translations of the unit length. That is, it is assumed that (Euclidean) length is translation invariant. The same is true for angular measurements, but now we have rotations instead of translations. And here comes the difference between angular and linear measurements. From the projective perspective translations along a line are projective transformations with one fixed point — a point at infinity, while rotations in a flat pencil of lines have no fixed lines.

This observation opens a way for a generalization [54]. To define a projective metric on a line, first of all we should select a projective transformation $A$ which will play the role of the unit translation. If $x$ is some point on a projective line then $x, A(x), A^2(x), A^3(x) \ldots$ will give mark-points of the distance scale. Therefore, we will have as many different measures of the distance on a projective line as there are substantially different projective transformations of this line. The Euclidean example suggests that we can classify projective transformations by the number of their fixed points.

On the projective line $P^1$ every three distinct points define the unique coordinate system in which the basic points have coordinates $(1,0)$, $(0,1)$ and $(1,1)$. Any other point is given by an equivalence class $[x] = (\lambda x_1, \lambda x_2)$, $\lambda \neq 0$, of two real numbers $x_1, x_2$. Projective transformations of the line $P^1$ have the form

$$\begin{align*}
x_1' &= A_{11}x_1 + A_{12}x_2, \\
x_2' &= A_{21}x_1 + A_{22}x_2 , \quad \text{det} (A_{ij}) \neq 0 .
\end{align*}$$

If $(x_1, x_2)$ is a fixed point of this transformation, then we should have

$$\begin{align*}
\lambda x_1 &= A_{11}x_1 + A_{12}x_2 , \\
\lambda x_2 &= A_{21}x_1 + A_{22}x_2 ,
\end{align*}$$

for some $\lambda \neq 0$. The point $(x_1, x_2)$ is uniquely determined by the ratio $z = \frac{x_1}{x_2}$ (the non-homogeneous coordinate of this point), for which we get from (24)

$$\begin{align*}
A_{21}z^2 + (A_{22} - A_{11})z - A_{12} &= 0 ,
\end{align*}$$
or in the homogeneous form
\[ \Omega(x, x) = 0 , \]  
(26)
where we have introduced the quadratic form
\[ \Omega(x, y) = \sum_{i,j=1}^{2} \Omega_{ij} x_i y_j , \]
with
\[ \Omega_{11} = A_{21} , \quad \Omega_{12} = \Omega_{21} = \frac{A_{22} - A_{11}}{2} , \quad \Omega_{22} = -A_{12} . \]
If \( \Delta = (A_{22} - A_{11})^2 + 4A_{12}A_{21} = 4(\Omega_{12}^2 - \Omega_{11}\Omega_{22}) > 0 \), then (25) has two different real solutions and the transformation \( A \) is called hyperbolic. If \( \Delta = 0 \), two different solutions degenerate into a single real solution and \( A \) is called parabolic. At last, if \( \Delta < 0 \), then (25) has no real solutions at all and the transformation \( A \) is called elliptic.

Therefore, we have three different measures of the distance on the projective line \( P^1 \): hyperbolic, parabolic and elliptic. By duality the same is true for a flat pencil of lines for which we have three different types of angular measure.

Let us consider first the hyperbolic measure. Then the corresponding projective transformation \( A \) has two fixed points \( p_0 \) and \( p_\infty \). Let these points be the basic points of the projective coordinate system such that \( p_0 = (0, 1) \) and \( p_\infty = (1, 0) \). Substituting these points into (24), we get \( A_{12} = A_{21} = 0 \) and, therefore, the projective transformation \( A \) takes the simple form
\[ x'_1 = A_{11} x_1 , \]
\[ x'_2 = A_{22} x_2 \]
in this coordinate system. Introducing again the non-homogeneous coordinate \( z = \frac{x_1}{x_2} \), we can represent the transformation \( A \) in the form \( z' = \lambda z \), where \( \lambda = A_{11}/A_{22} \).

If we take some point \( z = z_1 \) as the beginning of the distance scale then the mark-points of this scale will be \( z_1, \lambda z_1, \lambda^2 z_1, \lambda^3 z_1, \ldots \) and the distance between points \( \lambda^n z_1 \) and \( z_1 \) is equal to \( n \).

To measure distances that constitute fractions of the unit distance, one should subdivide the unit intervals \([\lambda^n z_1, \lambda^{n+1} z_1]\) into \( N \) equal parts. This can be achieved by means of the projective transformation \( z' = \mu z \), leaving the basic points \( p_0 (z = 0) \) and \( p_\infty (z = \infty) \) invariant, such that \( \mu^N z_1 = \lambda z_1 \), or \( \mu = \lambda^{\frac{1}{N}} \). Then the points \( z_1, \mu z_1, \mu^2 z_1, \ldots, \mu^N z_1 = \lambda z_1 \) constitute the desired subdivision of the unit interval \([z_1, \lambda z_1]\) into \( N \) equal parts. Now the distance from the point \( z_1 \) to the point \( \lambda^{n+\frac{m}{N}} z_1 \) is equal to \( n + \frac{m}{N} \).
Repeating subdivisions infinitely, we come to the conclusion [54] that the distance from $z_1$ to a point $z$ equals to the real number $\alpha$ such that $z = \lambda^\alpha z_1$. Therefore, in our particular coordinate system the hyperbolic distance is given by the formula

$$d(z, z_1) = \frac{1}{\ln \lambda} \ln \frac{z}{z_1}.$$  

As the distance on the line should be additive, $d(x, y) = d(x, z_1) + d(z_1, y)$, we get the general formula [54] for the hyperbolic distance between points $x = (x_1, x_2)$ and $y = (y_1, y_2)$

$$d(x, y) = C \ln \frac{x_1 y_2}{x_2 y_1},$$  \hspace{1cm} (27)

where $C = \frac{1}{\ln \lambda}$ is some constant and the freedom to choose $C$ reflects the freedom to choose different units of the length measurement.

However, as $p_0 = (0, 1)$ and $p_\infty = (1, 0)$, we have $x = x_2 p_0 + x_1 p_\infty$, $y = y_2 p_0 + y_1 p_\infty$ and

$$\frac{x_1 y_2}{x_2 y_1} = R(p_0, p_\infty, x, y)$$

is the cross-ratio formed by the points $x$ and $y$ with the fixed points $p_0$ and $p_\infty$ of the projective transformation $A$. Therefore we get the coordinate-independent form

$$d(x, y) = C \ln R(p_0, p_\infty, x, y)$$  \hspace{1cm} (28)

of the hyperbolic distance formula. Now we use this formula to express the hyperbolic distance in terms of homogeneous coordinates in a general coordinate system.

As the projective coordinates of a point are determined only up to a scale factor, we can write

$$p_0 = \lambda_0 x + y, \quad p_\infty = \lambda_\infty x + y,$$  \hspace{1cm} (29)

for some $\lambda_0$ and $\lambda_\infty$. Then, using [64] $R(p_0, p_\infty, x, y) = R(x, y, p_0, p_\infty)$, we get

$$R(p_0, p_\infty, x, y) = \frac{\lambda_\infty}{\lambda_0}.$$  

But $\Omega(p_0, p_0) = \Omega(p_\infty, p_\infty) = 0$ and (29) shows that both $\lambda_0$ and $\lambda_\infty$ are solutions of the quadratic equation

$$\lambda^2 \Omega(x, x) + 2\lambda \Omega(x, y) + \Omega(y, y) = 0.$$  \hspace{1cm} (30)
Therefore,
\[
R(p_0, p_\infty, x, y) = \frac{\Omega(x, y) + \sqrt{\Omega^2(x, y) - \Omega(x, x)\Omega(y, y)}}{\Omega(x, y) - \sqrt{\Omega^2(x, y) - \Omega(x, x)\Omega(y, y)}},
\]
and the hyperbolic distance is
\[
d(x, y) = C \ln \frac{\Omega(x, y) + \sqrt{\Omega^2(x, y) - \Omega(x, x)\Omega(y, y)}}{\Omega(x, y) - \sqrt{\Omega^2(x, y) - \Omega(x, x)\Omega(y, y)}}. \tag{31}
\]

We have not mentioned one subtlety [54]. The fixed points \(p_0\) and \(p_\infty\) divide the projective line \(P^1\) into two intervals according to the sign of the cross-ratio \(R(p_0, p_\infty, x, y)\). Above we have assumed that \(x\) and \(y\) are from the interval which corresponds to the positive sign of this cross-ratio. Therefore (31) gives the distance only between points of this interval. As this formula implies, \(p_0\) and \(p_\infty\) are both at logarithmically infinite distance from every point of the interval under consideration. Therefore these points, to say nothing of the points beyond them, are unreachable for the inhabitants of the one-dimensional hyperbolic world, for whom the question of existence of the other interval is a metaphysical question.

The analytic continuation of (31) by means of the formula
\[
\ln x = 2i \arccos \frac{x + 1}{2\sqrt{x}}
\]
can be used to get the distance in the elliptic case, when (30) does not have real solutions, and, therefore, \(\Omega^2(x, y) - \Omega(x, x)\Omega(y, y) < 0\) for all points \(x\) and \(y\). This gives
\[
d(x, y) = 2iC \arccos \frac{\Omega(x, y)}{\sqrt{\Omega(x, x)\Omega(y, y)}}. \tag{32}
\]

Note that in the elliptic case there are no infinite points, all distances being finite and defined only up to \(2\pi n, n \in \mathbb{Z}\). This is precisely the situation characteristic of the Euclidean angles. If we take \(\Omega(x, x) = x_1^2 + x_2^2\), so that \(\Omega_{ij} = \delta_{ij}\), then (32) gives for the distance (angle) between lines \((u_1, u_2)\) and \((v_1, v_2)\)
\[
d(u, v) = 2iC \arccos \frac{u_1 v_1 + u_2 v_2}{\sqrt{(u_1^2 + u_2^2)(v_1^2 + v_2^2)}}.
\]

Obviously this is the usual Euclidean formula for the angle between two lines with normals \((u_1, u_2)\) and \((v_1, v_2)\) provided that \(C = -i/2\).
It remains to consider the parabolic case, when $\Omega^2(x, y) - \Omega(x, x)\Omega(y, y)$ is identically zero for all $x, y$, and, therefore, (31) gives zero distance between every pair of points. Nevertheless it is possible to define a nontrivial parabolic distance if one considers how (31) approaches zero [54].

Noting that

$$\Omega(x, x)\Omega(y, y) - \Omega^2(x, y) = (x_1 y_2 - y_1 x_2) - \Delta,$$

where $\Delta = 4(\Omega_{12}^2 - \Omega_{11}\Omega_{22})$ is the discriminant of the quadratic form $\Omega$, we can rewrite (31) for small $\Delta$ as follows:

$$d(x, y) = 2iC \arcsin \sqrt{\frac{\Omega(x, x)\Omega(y, y) - \Omega^2(x, y)}{\Omega(x, x)\Omega(y, y)}} \approx iC \sqrt{-\Delta} \frac{x_1 y_2 - y_1 x_2}{\sqrt{\Omega(x, x)\Omega(y, y)}}.$$

However, we can assume that the arbitrary constant $C$ goes to infinity as $\Delta \to 0$ so that $iC \sqrt{-\Delta} = k$ remains finite and non zero. In the parabolic limit we have $\Omega(x, x) = (p_1 x_1 + p_2 x_2)^2$, where $p_1 = \sqrt{\Omega_{11}}$ and $p_2 = \sqrt{\Omega_{22}}$. Therefore, we get the following formula for the parabolic distance:

$$d(x, y) = k \frac{x_1 y_2 - y_1 x_2}{(p_1 x_1 + p_2 x_2)(p_1 y_1 + p_2 y_2)} = \frac{Q(x)}{P(x)} - \frac{Q(y)}{P(y)}, \quad (33)$$

where $P(x) = p_1 x_1 + p_2 x_2 = \sqrt{\Omega(x, x)}$ and $Q(x)$ is an arbitrary linear form not proportional to $P(x)$ (that is, $q_1 p_2 - q_2 p_1 = k \neq 0$). In particular, if $P(x) = x_2$ and $Q(x) = x_1$, we get the usual Euclidean expression for the distance between points $x$ and $y$ whose non-homogeneous coordinates are $\tilde{x} = \frac{x_1}{x_2}$ and $\tilde{y} = \frac{y_1}{y_2}$: $d(x, y) = \frac{x_1}{x_2} - \frac{y_1}{y_2} = \tilde{x} - \tilde{y}$.

7. Nine Cayley–Klein geometries

A point on the projective plane $P^2$ is determined by three real coordinates $x_1, x_2, x_3$. Therefore, a natural generalization of equation (26), which defines the fundamental points of the linear measure, is the equation of a conic section

$$\Omega(x, x) = \sum_{i,j=1}^{3} \Omega_{ij} x_i x_j = 0. \quad (34)$$

The conic section (34) will be called the fundamental conic, or the Absolute in Cayley’s terminology. Every line incident to an interior point of this conic section intersects it at two points, which can be used as the fundamental
points of a linear metric on the line. This allows us to define a projective metric for interior points of the Absolute. For a pair $x$ and $y$ of such points the line $l = x \times y$ intersects the Absolute at some points $p_0$ and $p_\infty$. As in the one-dimensional case we can write $p_0 = \lambda x + y$, $p_\infty = \lambda' x + y$ and applying the same reasoning we will end up with the same formula (31) for the hyperbolic distance between interior points of the Absolute lying on the line $l$. Of course, the arbitrary constant $C$ occurring in this formula must be the same for all lines on which (31) defines a metric.

Analogously, to define the angular measure, we can consider the Absolute as made from (tangential) lines

$$\tilde{\Omega}(u, u) = \sum_{i,j=1}^{3} \tilde{\Omega}_{ij} u_i u_j = 0,$$

(35)

where $(\tilde{\Omega}_{ij})$ is the matrix of cofactors [64] of the matrix $(\Omega_{ij})$. Then for every flat pencil of lines at every interior point of the Absolute we can choose two lines tangent to the Absolute (belonging to the Absolute considered as made from lines) as the fundamental lines to define the cross-ratio and the associated hyperbolic angular measure. If two fundamental lines degenerate into one line, we will have a parabolic angular measure and if the Absolute is such that the given pencil of lines does not contain (real) tangents to the Absolute, the angular measure will be elliptic. The arbitrary constant $C'$ which will appear in the projective angular measure must be the same for all pencils of lines but can be different from the constant $C$ which appears in the projective linear measure.

Therefore, we have nine different combinations of types of the linear and angular measures on the plane, and hence, nine different plane geometries. These geometries are nowadays called Cayley–Klein geometries and they are listed in the Table I.

**TABLE I**

<table>
<thead>
<tr>
<th>Measure of angles</th>
<th>Measure of lengths</th>
</tr>
</thead>
<tbody>
<tr>
<td>Hyperbolic</td>
<td>Hyperbolic</td>
</tr>
<tr>
<td>Doubly hyperbolic</td>
<td>Minkowski</td>
</tr>
<tr>
<td>(de-Sitter)</td>
<td>co-Hyperbolic</td>
</tr>
<tr>
<td></td>
<td>(anti de-Sitter)</td>
</tr>
<tr>
<td>Parabolic</td>
<td>co-Minkowski</td>
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<tr>
<td></td>
<td>Galilean</td>
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<tr>
<td></td>
<td>co-Euclidean</td>
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<tr>
<td>Elliptic</td>
<td>Hyperbolic</td>
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<tr>
<td>(Lobachevsky)</td>
<td>Euclidean</td>
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<tr>
<td></td>
<td>Elliptic</td>
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<tr>
<td></td>
<td>(Riemann)</td>
</tr>
</tbody>
</table>
Let us take a closer look at these geometries and at relations among them. First of all, some geometries are related by duality. The points of geometry \(X\) are the lines of its dual, or co-geometry \(\tilde{X}\), and the lines of geometry \(X\) are the points of \(\tilde{X}\). What was the distance between points in geometry \(X\), in the dual geometry \(\tilde{X}\) we should call the angle between lines, and what was called the angle between lines in \(X\), becomes in \(\tilde{X}\) the distance between points. Elliptic, Galilean and doubly hyperbolic geometries are self-dual, while all other geometries differ from their duals.

There exists [64] a projective coordinate system in which the equation of the conic \(\Omega\) takes the simple form

\[
\sum_{i=1}^{3} b_i x_i^2 = 0, \quad \text{or in the tangential form } \sum_{i=1}^{3} \frac{u_i^2}{b_i} = 0.
\]

The homogeneous coordinates are determined only up to a scale factor. Therefore, we can assume without loss of generality the following equation [68] for the conic \(\Omega\):

\[
\Omega(x, x) = \epsilon_1 \epsilon_2 x_1^2 - \epsilon_1 x_2^2 + x_3^2 = 0,
\]

(36)

with some constants \(\epsilon_1, \epsilon_2\), or in the tangential form,

\[
\tilde{\Omega}(l, l) = l_1^2 - \epsilon_2 l_2^2 + \epsilon_1 \epsilon_2 l_3^2 = 0.
\]

(37)

Let us introduce the following non-homogeneous coordinates (in the Euclidean case they will turn out to be the usual Cartesian coordinates)

\[
z = \frac{x_1}{\sqrt{\Omega(x, x)} + x_3}, \quad t = \frac{x_2}{\sqrt{\Omega(x, x)} + x_3}.
\]

(38)

If the quadratic form \(\Omega\) is not positive definite, we will assume that only those points for which \(\Omega(x, x) > 0\) are the points of the corresponding geometry. however, if \(\Omega(x, x) > 0\) then we can choose an arbitrary scale factor of the homogeneous coordinates so that

\[
\Omega(x, x) = 1
\]

(39)

for the scaled homogeneous coordinates. From (38) and (39) we get expressions of the corresponding values of the homogeneous coordinates in terms of \(z\) and \(t\):

\[
x_1 = \frac{2z}{1 - \epsilon_1 (t^2 - \epsilon_2 z^2)}, \quad x_2 = \frac{2t}{1 - \epsilon_1 (t^2 - \epsilon_2 z^2)}, \quad x_3 = \frac{1 + \epsilon_1 (t^2 - \epsilon_2 z^2)}{1 - \epsilon_1 (t^2 - \epsilon_2 z^2)}.
\]

(40)
Taking $C = k/(2\sqrt{\epsilon_1})$, we get for the distance between points $x_0(t_0, z_0)$ and $x(t, z)$ from (32)
\[
\cosh \left( \frac{\sqrt{\epsilon_1}}{k} d(x_0, x) \right) = \Omega(x_0, x),
\]
which by using $\sinh^2 (x/2) = (\cosh x - 1)/2$ can be transformed into
\[
\sinh^2 \left( \frac{\sqrt{\epsilon_1}}{2k} d(x_0, x) \right) = \epsilon_1 \frac{(t - t_0)^2 - \epsilon_2(z - z_0)^2}{[1 - \epsilon_1(t^2 - \epsilon_2 z^2)][1 - \epsilon_1(t_0^2 - \epsilon_2 z_0^2)]}. \tag{41}
\]

It is well known that the points of the Euclidean plane can be represented by complex numbers. Remarkably, the points of all nine Cayley–Klein geometries also can be represented by suitably generalized complex numbers [69, 70].

Complex numbers $a + ib$ are obtained by adding a special element $i$ to the real numbers. This special element is characterized by the property that it is a solution of the quadratic equation $i^2 = -1$. However, there is nothing special to this quadratic equation. On equal footing we can assume a special element $e$ to be a solution of the general quadratic equation $Ae^2 + Be + C = 0$. In fact this construction gives three different types of generalized complex numbers $a + eb$ depending on the value of the discriminant: $\Delta = B^2 - 4AC < 0$, $\Delta = 0$, or $\Delta > 0$. In the first case we get the ordinary complex numbers $a + ib$ and one can assume without loss of generality that $i^2 = -1$. If the discriminant is zero, we get the so called dual numbers $a + \epsilon b$ and one can assume that $\epsilon^2 = 0$. At last, for a positive discriminant we get the double numbers $a + eb$ with $e^2 = 1$. Note that, for example, in the double numbers $e$ is a special unit different from 1 or $-1$. That is, the equation $x^2 = 1$ has four different solutions in double numbers. All that is explained in detail in Yaglom’s book *Complex Numbers in Geometry* [69].

Now let us introduce a special element $e$ such that $e^2 = \epsilon_2$ and the corresponding generalized complex numbers $z = t + ex$. Then (41) yields for the distance between points $z_0 = t_0 + ex_0$ and $z$
\[
\sinh^2 \left( \frac{\sqrt{\epsilon_1}}{2k} d(z_0, z) \right) = \epsilon_1 \frac{(z - z_0)(\bar{z} - \bar{z}_0)}{[1 - \epsilon_1(z\bar{z})][1 - \epsilon_1(z_0\bar{z}_0)]}, \tag{42}
\]
where the conjugation operation is defined as usual: $\bar{z} = t - ex$ and $z\bar{z} = t^2 - e^2 x^2 = t^2 - \epsilon_2 x^2$.

Analogous considerations apply to the angular measure based on the tangential conic (37) and we get similar formulas for the angle between two lines, with the roles of $\epsilon_1$ and $\epsilon_2$ interchanged.

In fact for $\epsilon_1$ there are only three possibilities: $\epsilon_1 = 1$, $\epsilon_1 = 0$ or $\epsilon_1 = -1$. All other cases can be reduced to these three by changing the unit of the
linear measure. The same is true for the parameter $\epsilon_2$ also: by changing the unit of the angular measure this parameter can be brought to 1 or $-1$, if different from zero.

If $\epsilon_1 = 1$, then (41) shows that the linear measure is hyperbolic, the points of the corresponding Cayley–Klein geometry are represented by double numbers if $\epsilon_2 = 1$ (hyperbolic angular measure, the de Sitter geometry), by dual numbers if $\epsilon_2 = 0$ (parabolic angular measure, the co-Minkowski geometry), and by complex numbers if $\epsilon_2 = -1$ (elliptic angular measure, the Lobachevsky geometry). The corresponding distance formula is

$$\sinh^2 \frac{d(z_0, z)}{2k} = \frac{(\bar{z} - \bar{z}_0)(\bar{z} - \bar{z}_0)}{[1 - \bar{z}\bar{z}][1 - \bar{z}_0\bar{z}_0]}.$$  

If $\epsilon_1 = -1$ then we have elliptic linear measure and the distance formula is

$$\sin^2 \frac{d(z_0, z)}{2k} = \frac{(\bar{z} - \bar{z}_0)(\bar{z} - \bar{z}_0)}{[1 + \bar{z}\bar{z}][1 + \bar{z}_0\bar{z}_0]}.$$  

At that, the points of the anti de-Sitter geometry are represented by double numbers, points of the co-Euclidean geometry — by dual numbers, and points of the Elliptic geometry — by complex numbers. For $\epsilon_1 = \pm 1$, the choice $k = 1$ corresponds to the usual definition of length in the Lobachevsky and Elliptic (Riemann) geometries [70].

At last, taking the limit $\epsilon_1 \to 0$, we get the distance in the parabolic case

$$\frac{d^2(z_0, z)}{4k^2} = (\bar{z} - \bar{z}_0)(\bar{z} - \bar{z}_0).$$  

In the case of the elliptic angular measure, $z$ and $z_0$ are complex numbers, and the last formula gives the usual Euclidean distance if we take $k = 1/2$. The points of the Galilean geometry (parabolic angular measure) are represented by dual numbers and the points of the Minkowski geometry (hyperbolic angular measure) — by double numbers [70].

We can make contact with Riemannian geometry by noting that for infinitesimally close $z$ and $z_0$, and for the unit scale factor $k = 1$, (41) takes the form

$$ds^2 = \frac{4(dt^2 - \epsilon_2 dx^2)}{[1 - \epsilon_1(t^2 - \epsilon_2 x^2)]^2}. \quad (43)$$  

Let us define the generalized cosine and sine functions as follows [71, 72]

$$C(x; \epsilon) = \begin{cases} \cos (\sqrt{-\epsilon} x), & \text{if } \epsilon < 0 \\ 1, & \text{if } \epsilon = 0 \\ \cosh (\sqrt{\epsilon} x) & \text{if } \epsilon > 0 \end{cases}, \quad S(x; \epsilon) = \begin{cases} \frac{\sin (\sqrt{-\epsilon} x)}{\sqrt{-\epsilon}}, & \text{if } \epsilon < 0 \\ x, & \text{if } \epsilon = 0 \\ \frac{\sinh (\sqrt{\epsilon} x)}{\sqrt{\epsilon}}, & \text{if } \epsilon > 0 \end{cases}. \quad (44)$$
Then
\[ C^2(x; \epsilon) - \epsilon S^2(x; \epsilon) = 1, \quad \frac{dC(x; \epsilon)}{dx} = \epsilon S(x; \epsilon), \quad \frac{dS(x; \epsilon)}{dx} = C(x; \epsilon). \quad (45) \]

If we define the “polar” coordinates \((r, \phi)\) (for points with \(0 < t^2 - \epsilon_2 x^2 < \frac{1}{\epsilon_1}\)) through relations
\[ t = rC(\phi; \epsilon_2), \quad x = rS(\phi; \epsilon_2), \]
then, by using identities (45), we obtain from (43)
\[ ds^2 = \frac{4(dr^2 - \epsilon_2 r^2 d\phi^2)}{[1 - \epsilon_1 r^2]^2} = E(r)dr^2 + G(r)d\phi^2, \quad (46) \]
where
\[ E(r) = \frac{4}{[1 - \epsilon_1 r^2]^2} \quad \text{and} \quad G(r) = -\frac{4\epsilon_2 r^2}{[1 - \epsilon_1 r^2]^2}. \quad (47) \]

However, for the Riemannian metric (46) the corresponding Gaussian curvature \(K\) satisfies the equation [54]
\[ 4E^2G^2K = E \left( \frac{dG}{dr} \right)^2 + Gin \frac{dE}{dr} \frac{dG}{dr} - 2EG \frac{d^2G}{dr^2}. \]

Substituting (47), we get
\[ K = -\epsilon_1. \]

This clarifies the geometric meaning of the parameters \(\epsilon_1\) and (by duality) \(\epsilon_2\): the corresponding Cayley–Klein geometry has the constant curvature \(-\epsilon_1\) and its dual geometry, the constant curvature \(-\epsilon_2\).

More details about the Cayley–Klein geometries can be found in [70], or from more abstract group-theoretical point of view, in [71–74]. We will not follow this abstract trend here, but for the sake of future use give below the derivation of the Lie algebras of the Cayley–Klein symmetry groups.

Rigid motions (symmetries) of the Cayley–Klein geometry are those projective transformations of \(P^2\) which leave the fundamental conic (36) invariant. Let \(S = e^{\alpha G}\) be such a transformation with \(G\) as its generator (an element of the Lie algebra of the Cayley–Klein geometry symmetry group). Writing (36) in the matrix form \(x^T \Omega x = 0\), for rigid motions we get \((x')^T \Omega x' = x^T \Omega x\), where \(x' = Sx\). Therefore, the invariance of the Absolute \(\Omega\) under the symmetry transformation \(S\) implies that \(S^T \Omega S = \Omega\). For the infinitesimal parameter \(\alpha\) this condition reduces to
\[ G^T \Omega + \Omega G = 0. \]
Taking in this equation

\[ \Omega = \begin{pmatrix} \epsilon_1 \epsilon_2 & 0 & 0 \\ 0 & -\epsilon_1 & 0 \\ 0 & 0 & 1 \end{pmatrix}, \]

we get (in the general case \( \epsilon_1 \epsilon_2 \neq 0 \))

\[ G_{11} = G_{22} = G_{33} = 0, \quad G_{21} = \epsilon_2 G_{12}, \quad G_{31} = -\epsilon_1 \epsilon_2 G_{13}, \quad G_{32} = \epsilon_1 G_{23}. \]

Therefore, the Lie algebra of the symmetry group of the Cayley–Klein geometry has three linearly independent generators [68]

\[ G_1 = \begin{pmatrix} 0 & 0 & 0 \\ 0 & 0 & -1 \\ 0 & -\epsilon_1 & 0 \end{pmatrix}, \quad G_2 = \begin{pmatrix} 0 & 0 & 1 \\ 0 & 0 & 0 \\ -\epsilon_1 \epsilon_2 & 0 & 0 \end{pmatrix}, \]

\[ G_3 = \begin{pmatrix} 0 & 1 & 0 \\ \epsilon_2 & 0 & 0 \\ 0 & 0 & 0 \end{pmatrix}. \quad (48) \]

Now it is not difficult to find the commutators

\[ [G_1, G_2] = \epsilon_1 G_3, \quad [G_3, G_1] = -G_2, \quad [G_2, G_3] = \epsilon_2 G_1. \quad (49) \]

In fact all relevant information is encoded in these commutation relations. Let us consider, for example, the Minkowski geometry with \( \epsilon_1 = 0, \epsilon_2 = 1 \) and obtain the Lorentz transformations from its Lie algebra.

For the Minkowski geometry, the commuting generators \( G_1 = H \) and \( G_2 = P \) can be considered as the generators of time and space translations, while \( G_3 = K \) will play the role of the Lorentz transformation generator. From (49) we read the commutation relations

\[ [H, P] = 0, \quad [K, H] = -P, \quad [K, P] = -H. \quad (50) \]

among these generators. This defines the Lie algebra of the two-dimensional Poincaré group \( \mathcal{P} \). The two-dimensional Lorentz group \( \mathcal{L} \) is generated by \( K \) and consists of the transformations of the form \( e^{\psi K} \), where \( \psi \) is the relevant group parameter. As we stated above, the Minkowski space-time can be identified with the coset space \( \mathcal{M} = \mathcal{P}/\mathcal{L} \). Every element of \( \mathcal{M} \) has the form \([g] = g \mathcal{L} \), where \( g = e^{x_0 H} e^{x P} = e^{x_0 H + x P} \) is the element of \( \mathcal{P} \) characterized by two real parameters \( x_0 \) and \( x \), which we identify with the time and space coordinates of the “point” \([g]\). The Lorentz transformation \( s = e^{\psi K} \) acts on the point \([g]\) in the following way

\[ [g] \rightarrow [sg] = [sgs^{-1}], \]
where the last equality follows from the fact that \( s^{-1} \in \mathcal{L} \), and therefore, 
\( s^{-1} \mathcal{L} = \mathcal{L} \). Now we need to calculate

\[
e^{\psi K} e^{x_0 H + xP} e^{-\psi K}.
\]

First of all, note that (48) implies that \( H^2 = P^2 = HP = 0 \), and therefore,

\[
e^{x_0 H + xP} = 1 + x_0 H + xP.
\]

To proceed, recall the Baker–Campbell–Hausdorff formula [75]

\[e^A B e^{-A} = \sum_{m=0}^{\infty} \frac{B_m}{m!}, \quad (51)\]

where \( B_m \) is defined recursively by \( B_m = [A, B_{m-1}] \) and \( B_0 = B \). Taking \( A = \psi K, B = x_0 H + xP \), we get from (50)

\[B_1 = [A, B] = -\psi (x_0 P + xH), \quad B_2 = [A, B_1] = \psi^2 (x_0 H + xP) = \psi^2 B.
\]

Therefore,

\[B_{2k+1} = -\psi^{2k+1} (x_0 P + xH), \quad B_{2k} = \psi^{2k} (x_0 H + xP)
\]

and

\[
\sum_{m=0}^{\infty} \frac{B_m}{m!} = \sum_{k=0}^{\infty} \left\{ \frac{\psi^{2k}}{(2k)!} (x_0 H + xP) - \frac{\psi^{2k+1}}{(2k+1)!} (x_0 P + xH) \right\}.
\]

Or, after summing up the infinite series,

\[
\sum_{m=0}^{\infty} \frac{B_m}{m!} = (\cosh \psi x_0 - \sinh \psi x) H + (\cosh \psi x - \sinh \psi x_0) P.
\]

Hence, if we write \( sgs^{-1} = e^{x'_0 H} e^{x'P} = 1 + x'_0 H + x'P \), that is, if we represent the point \([sg]\) by the transformed coordinates \( x'_0, x' \), we get

\[
x'_0 = \cosh \psi x_0 - \sinh \psi x,
\]

\[
x' = \cosh \psi x - \sinh \psi x_0,
\]

which is nothing but the Lorentz transformation (13) provided \( \psi \) is the rapidity defined by \( \tanh \psi = \beta \), and \( x_0 = ct \).

This derivation of Lorentz transformations demonstrates clearly that the natural (canonical) parameter associated with Lorentz transformations is rapidity, not velocity [76–78]. Therefore, to follow the intrinsic instead of historical logic, we should “introduce rapidity as soon as possible in the teaching of relativity, namely, at the start. There is no need to go through the expressions of Lorentz transformations using velocity, and then to ‘discover’ the elegant properties of rapidity as if they resulted from some happy and unpredictable circumstance” [76].
8. Possible kinematics

As we see, special relativity has its geometric roots in the Cayley–Klein geometries. However, special relativity is not a geometric theory, but a physical one. This means that it includes concepts like causality, reference frames, inertial motion, relativity principle, which, being basic physical concepts, are foreign to geometry. We showed in the first chapters that if we stick to these concepts then the Minkowski geometry and its singular $c \to \infty$ Galilean cousin remain as the only possibilities.

But what about the other Cayley–Klein geometries? Do they also have physical meanings? To answer this question affirmatively we have to alter or modify the physical premises of special relativity and, as the preceding brief discussion of relativity without reference frames indicates, we should first of all generalize the Relativity Principle, getting rid of the too restrictive framework of inertial reference frames, whose existence in general space-times is neither obvious nor guaranteed.

The symmetry group of special relativity is the ten-parameter Poincaré group. Ten basis elements of its Lie algebra are the following: the generator $H$ of time translations; three generators $P_i$ of space translations along the $i$-axis; three generators $J_i$ of spatial rotations; and three generators $K_i$ of pure Lorentz transformations, which can be considered as the inertial transformations (boosts) along the $i$-axis. The commutation relations involving $J_i$ have the form

$$[J_i, H] = 0, \quad [J_i, J_j] = \epsilon_{ijk} J_k, \quad [J_i, P_j] = \epsilon_{ijk} P_k, \quad [J_i, K_j] = \epsilon_{ijk} K_k,$$

and they just tell us that $H$ is a scalar and $P_i, J_i, K_i$ are vectors. These commutation relations can not be altered without spoiling the isotropy of space if we still want to regard $H$ as a scalar and $P_i, J_i, K_i$ as vectors. However, other commutation relations

$$[H, P_i] = 0, \quad [H, K_i] = P_i, \quad [P_i, P_j] = 0, \quad [K_i, K_j] = -\epsilon_{ijk} J_k, \quad [P_i, K_j] = \delta_{ij} H,$$

are less rigid as they depend on the interpretation of inertial transformations which we want to change.

Besides the continuous symmetries of the Poincaré group there are some discrete symmetries such as the space inversion (parity) and the time-reversal, which play important roles in physics. Under the time-reversal

$$H \mapsto -H, \quad P \mapsto P$$

and the commutation relations (53) and (54) indicate that (55) can be extended up to the involutive automorphism

$$\tau: \quad H \mapsto -H, \quad P \mapsto P, \quad J \mapsto J, \quad K \mapsto -K$$
of the Poincaré Lie Algebra. Analogously, the space inversion is represented by the involutive automorphism

$$\pi : \ H \mapsto H, \ P \mapsto -P, \ J \mapsto J, \ K \mapsto -K. \quad (57)$$

We will assume [79] that the generalized Lie algebra we are looking for also possesses the automorphisms $\tau$ and $\pi$. From the physical point of view we are assuming that the observed $P$- and $T$-asymmetries of the weak interactions have no geometric origin and the space-time itself is mirror symmetric. Of course, it is tempting then to expect the world of elementary particles to be also mirror symmetric and the less known fact is that this mirror symmetry can be indeed ensured by introducing hypothetical mirror matter counterparts of the ordinary elementary particles [80].

After deforming the commutation relations (54), the Poincaré group is replaced by the so called kinematical group — the generalized relativity group of nature. In a remarkable paper [79] Bacry and Lévy–Leblond showed that under very general assumptions there are only eleven possible kinematics. The assumptions they used are:

- The infinitesimal generators $H, P_i, J_i, K_i$ transform correctly under rotations as required by the space isotropy. This leads to the commutation relations (53).

- the parity $\pi$ and the time-reversal $\tau$ are automorphisms of the kinematical group.

- Inertial transformations in any given direction form a noncompact subgroup. Otherwise, a sufficiently large boost would be no boost at all, like $4\pi$ rotations [81], contrary to the physical meaning we ascribe to boosts. The role of this condition is to exclude space-times with no causal order, that is, timeless universes like the Euclidean one. Although, as was mentioned earlier, the real space-time quite might have such inclusions, but as there is no time in these regions and hence no motion, it is logical that the corresponding symmetry groups are not called kinematical.

If we demand the parity and time-reversal invariance, the only possible deformations of the commutation relations (54) will have the form

$$[H, P_i] = \epsilon_1 K_i, \quad [H, K_i] = \lambda P_i, \quad [P_i, P_j] = \alpha \epsilon_{ijk} J_k, \quad [K_i, K_j] = \beta \epsilon_{ijk} J_k, \quad [P_i, K_j] = \epsilon_2 \delta_{ij} H. \quad (58)$$
Now we have to ensure the Jacobi identities
\[
[P_i, [P_j, K_k]] + [P_j, [K_k, P_i]] + [K_k, [P_i, P_j]] = 0
\]
and
\[
[P_i, [K_j, K_k]] + [K_j, [K_k, P_i]] + [K_k, [P_i, K_j]] = 0
\]
which are satisfied only if
\[
\alpha - \epsilon_1 \epsilon_2 = 0 \tag{59}
\]
and
\[
\beta + \lambda \epsilon_2 = 0 \tag{60}
\]
It turns out [79] (this can be checked by explicit calculations) that all other Jacobi identities are also satisfied if (59) and (60) hold.

As we see, the structure of the generalized Lie algebra is completely determined by three real parameters \( \epsilon_1, \epsilon_2 \) and \( \lambda \). Note that the overall sign of the structure constants is irrelevant as the sign change of all structure constants can be achieved simply by multiplying each infinitesimal generator by \(-1\). Therefore we can assume \( \lambda \geq 0 \) without loss of generality and by a scale change it can be brought either to \( \lambda = 1 \) or \( \lambda = 0 \).

If \( \lambda = 1 \), the commutation relations are
\[
[H, P_i] = \epsilon_1 K_i, \quad [H, K_i] = P_i, \quad [P_i, P_j] = \epsilon_1 \epsilon_2 \epsilon_{ijk} J_k,
\]
\[
[K_i, K_j] = -\epsilon_2 \epsilon_{ijk} J_k, \quad [P_i, K_j] = \epsilon_2 \delta_{ij} H
\]
and comparing with (49) we see that every three-dimensional subalgebra \( (H, P_1, K_1) \), \( (H, P_2, K_2) \) and \( (H, P_3, K_3) \) realizes the Cayley–Klein geometry of the type \( (\epsilon_1, \epsilon_2) \) if the identifications \( G_1 = H, G_2 = P_i, G_3 = K_i \) are made.

Let us note that \( G_3^2 = \epsilon_2 E_{12} \) and \( G_3 E_{12} = E_{12} G_3 = G_3 \), where
\[
E_{12} = \begin{pmatrix} 1 & 0 & 0 \\ 0 & 1 & 0 \\ 0 & 0 & 0 \end{pmatrix}.
\]
Therefore,
\[
e^{\psi K} = 1 + \left( \psi + \frac{\psi^3}{3!} \epsilon_2 + \frac{\psi^5}{5!} \epsilon_2^2 + \frac{\psi^7}{7!} \epsilon_2^3 + \ldots \right) K
\]
\[
+ \left( \frac{\psi^2}{2!} \epsilon_2 + \frac{\psi^4}{4!} \epsilon_2^2 + \frac{\psi^6}{6!} \epsilon_2^3 + \ldots \right) E_{12}
\]
and recalling the definition of the generalized sine and cosine functions (44) we get
\[ e^{\psi K} = (1 - E_{12}) + S(\psi; \epsilon_2)K + C(\psi; \epsilon_2)E_{12}. \] (62)
As we see, \( e^{\psi K} \) is periodic if \( \epsilon_2 < 0 \); that is, in this case the inertial transformations form a compact group and the corresponding Cayley–Klein geometries with elliptic angular measure (Lobachevsky, Euclidean and Elliptic) do not lead to kinematical groups.

The remaining six cases include: the de Sitter kinematics (DS) with doubly-hyperbolic geometry \( \epsilon_1 > 0, \epsilon_2 > 0 \), the anti de Sitter kinematics (ADS) with the co-hyperbolic geometry \( \epsilon_1 < 0, \epsilon_2 > 0 \), the Poincaré kinematics \( (P) \) with the Minkowski geometry \( \epsilon_1 = 0, \epsilon_2 > 0 \), the Newton–Hook kinematics \( (NH) \) with the co-Minkowski geometry \( \epsilon_1 > 0, \epsilon_2 = 0 \), the anti Newton–Hook kinematics \( (ANH) \) with the co-Euclidean geometry \( \epsilon_1 < 0, \epsilon_2 = 0 \) and the Galilean kinematics \( (G) \) \( \epsilon_1 = 0, \epsilon_2 = 0 \) whose geometry is described in detail in [70].

At last, if \( \lambda = 0 \), the commutation relations are
\[ [H, P_i] = \epsilon_1 K_i, \quad [H, K_i] = 0, \quad [P_i, P_j] = \epsilon_1 \epsilon_2 \epsilon_{ijk} J_k, \]
\[ [K_i, K_j] = 0, \quad [P_i, K_j] = \epsilon_2 \delta_{ij} H. \] (63)
Suppose that \( \epsilon_1 \neq 0 \), and introduce another basis in the Lie algebra (63):
\[ P'_i = \epsilon_1 K_i, \quad K'_i = P_i, \quad H' = H, \quad J'_i = J_i. \]
The commutation relations in the new bases take the form
\[ [H', P'_i] = 0, \quad [H', K'_i] = P', \quad [P'_i, P'_j] = 0, \]
\[ [K'_i, K'_j] = \epsilon_1 \epsilon_2 \epsilon_{ijk} J'_k, \quad [P'_i, K'_j] = -\epsilon_1 \epsilon_2 \delta_{ij} H'. \]
However, this is the same Lie algebra as (61) for the Cayley–Klein parameters \( \epsilon'_1 = 0 \) and \( \epsilon'_2 = -\epsilon_1 \epsilon_2 \). Nevertheless, the physics corresponding to the isomorphic algebras (63) and (61) are completely different because we prescribe to the generators \( H, P, K \) a well-defined concrete physical meaning and they cannot be transformed arbitrarily, except by scale changes [79].

Anyway, all these new possibilities also realize Cayley–Klein geometries. At that, \( K_i = G_2/\epsilon_1 \) and by using \( G_2^2 = 0 \) (for \( \epsilon'_1 = 0 \)) we get \( e^{\psi K_i} = 1 + \psi K_i \); that is, the subgroup generated by \( K_i \) is always non-compact in this case. However, \( (-\epsilon_1, -\epsilon_2) \) and \( (\epsilon_1, \epsilon_2) \) give the same Lie algebra as the first case just corresponds to the change of basis: \( H \mapsto -H, \ P \mapsto -P, \ K \mapsto -K, \ J \mapsto J \). Therefore, one is left with three possibilities: the anti para-Poincaré kinematics \( (AP') \) \( \epsilon_1 = 1, \epsilon_2 = 1 \) with Euclidean geometry \( \epsilon'_1 = 0, \epsilon'_2 = -1 \), the para-Poincaré kinematics \( (P') \) \( \epsilon_1 = -1, \epsilon_2 = 1 \) with the Minkowski
geometry $\epsilon_1 = 0$, $\epsilon_2 = 1$, and the para-Galilei kinematics ($G'$) $\epsilon_1 = 1$, $\epsilon_2 = 0$ with Galilean geometry $\epsilon_1' = 0$, $\epsilon_2' = 0$.

The case $\lambda = 0$, $\epsilon_1 = 0$ adds two more possibilities: the Carroll kinematics ($C$) with $\epsilon_2 = \pm 1$, first discovered in [82], and the static kinematics ($S$) with $\epsilon_2 = 0$ for which all commutators between $H$, $P$, $K$ vanish.

In the case of the Carroll kinematics, the Galilean geometry is realized in each $(H, P_i, K_i)$ subspaces with $G_1 = P_i$, $G_2 = H$, $G_3 = K_i$. As we see, compared to the Galilean kinematics, the roles of time and space translations are interchanged in the Carroll kinematics. This leads to the exotic situation of absolute space but relative time. That is, an event has the same spatial coordinates irrespective of the applied inertial transformation (change of the reference frame) [82, 83]. There are no interactions between spatially separated events, no true motion, and practically no causality. The evolution of isolated and immobile physical objects corresponds to the ultralocal approximation of strong gravity [84]. The name of this strange kinematics is after Lewis Carroll’s tale Through the Looking-Glass, and What Alice Found There (1871), where the Red Queen points out to Alice: “A slow sort of country! Now, here, you see, it takes all the running you can do, to keep in the same place. If you want to get somewhere else, you must run at least twice as fast as that!”

It is remarkable that under very general assumptions all possible kinematics, except the static one, have Cayley–Klein geometries at their background [85, 86]. By deforming the Absolute of the doubly-hyperbolic geometry (by changing parameters $\epsilon_1$ and $\epsilon_2$) one can get all other Cayley–Klein geometries. It is not surprising, therefore, that all eleven possible kinematics are in fact limiting cases of the de Sitter or anti de Sitter kinematics [79] which are the most general relativity groups. As a result, there exists essentially only one way to generalize special relativity, namely, by endowing space-time with some constant curvature [79].

9. Group contractions

There is a subtlety in considering various limiting cases of the kinematical groups introduced in the previous chapter. We demonstrate this by considering the non-relativistic limit of the Poincaré group. It is commonplace that the Lorentz transformations (13), written for the space and time intervals $\Delta x$ and $\Delta t$,

$$
\Delta x' = \frac{1}{\sqrt{1 - \beta^2}}(\Delta x - V \Delta t),
$$

$$
\Delta t' = \frac{1}{\sqrt{1 - \beta^2}} \left(\Delta t - \frac{V}{c^2} \Delta x\right),
$$
in the non-relativistic limit $\beta \to 0$ reduce to the Galilei transformations
\begin{equation}
\Delta x' = \Delta x - V \Delta t,
\Delta t' = \Delta t.
\end{equation}

But to reach this conclusion, besides $\beta \ll 1$, we have implicitly assumed that
\begin{align*}
\beta \frac{\Delta x}{c \Delta t} & \ll 1, & \beta \frac{c \Delta t}{\Delta x} & \sim 1,
\text{or}
\beta & \ll \frac{c \Delta t}{\Delta x}, & \frac{\Delta x}{c \Delta t} & \sim 1,
\end{align*}
which is not necessarily true if $\Delta x \gg c \Delta t$.

Therefore the Galilei transformations are not the non-relativistic limit of the Lorentz transformations but only one non-relativistic limit, which corresponds to the situation when space intervals are much smaller than time intervals [82] (in units where $c = 1$). There exists another non-relativistic limit
\begin{align*}
\beta \frac{c \Delta t}{\Delta x} & \ll 1, & \beta \frac{\Delta x}{c \Delta t} & \sim 1,
\text{or}
\beta & \ll \frac{\Delta x}{c \Delta t}, & \frac{c \Delta t}{\Delta x} & \sim \beta \ll 1
\end{align*}
and in this limit the Lorentz transformations reduce to the Carroll transformations [82, 83]
\begin{align*}
\Delta x' &= \Delta x, \\
\Delta t' &= \Delta t - V c^2 \Delta x.
\end{align*}

The Carroll kinematics corresponds to the situation when space intervals are much larger than time intervals, hence no causal relationship between events are possible; events are isolated.

A systematic way to correctly treat limiting cases of symmetry groups was given by Inönu and Wigner [87] and is called “group contraction”. A symmetry group $G$ can be contracted towards its continuous subgroup $S$, which remain intact under the contraction process. Denote by $J_i$ the generators of the subgroup $S$ and the remaining generators of $G$, by $I_i$. Therefore the Lie algebra of the group $G$ looks like
\begin{align*}
[J_i, J_j] &= f^{(1)}_{ijk} J_k, & [I_i, J_j] &= f^{(2)}_{ijk} J_k + g^{(2)}_{ijk} I_k, & [I_i, I_j] &= f^{(3)}_{ijk} J_k + g^{(3)}_{ijk} I_k.
\end{align*}
Under the change of basis
\[ J'_i = J_i, \quad I'_i = \epsilon I_i, \] (67)
the commutation relations transform to
\[
[J'_i, J'_j] = f^{(1)}_{ijk} J'_k, \quad [I'_i, J'_j] = \epsilon f^{(2)}_{ijk} J'_k + g^{(2)}_{ijk} I'_k, \quad [I'_i, I'_j] = \epsilon^2 f^{(3)}_{ijk} J'_k + \epsilon g^{(3)}_{ijk} I'_k.
\] (68)

When \( \epsilon \to 0 \), the base change (67) becomes singular but the commutation relations (68) still have a well-defined limit as \( \epsilon \to 0 \):
\[
[J'_i, J'_j] = f^{(1)}_{ijk} J'_k, \quad [I'_i, J'_j] = g^{(2)}_{ijk} I'_k, \quad [I'_i, I'_j] = 0.
\] (69)

In general, the Lie algebra (69) is not isomorphic to the initial Lie algebra (66) and defines another symmetry group \( G' \), which is said to be the result of the contraction of the group \( G \) towards its continuous subgroup \( S \). The contracted generators \( I'_i \) form an Abelian invariant subgroup in \( G' \), as (69) shows.

The contraction has a clear meaning in the group parameter space. If \( I' = \epsilon I \), then the corresponding group parameters should satisfy \( \alpha = \epsilon \alpha' \) if we want \( e^{\alpha I} \) and \( e^{\epsilon \alpha' I'} \) to represent the same point (transformation) of the group. Therefore, when \( \epsilon \to 0 \) the parameter \( \alpha \) becomes infinitesimal. From the point of view of \( G \), its contracted form \( G' \) covers only infinitesimally small neighborhood of the subgroup \( S \). This explains why the operation is called contraction.

Let us return to kinematical groups. We do not want to spoil the space isotropy by contraction. Therefore, \( S \) should be a rotation-invariant subgroup of \( G \). Looking at the commutation relations (53) and (58), we find only four rotation-invariant subalgebras, generated respectively by \((J_i, H)\), \((J_i, P_i)\), \((J_i, K_i)\) and \((J_i)\), which are common to all kinematical Lie algebras. Therefore, we can consider four types of physical contractions of the kinematical groups [79]:

- Contraction with respect to the rotation and time-translation subgroups generated by \((J_i, H)\). Under this contraction
  \[ P_i \to \epsilon P_i, \quad K_i \to \epsilon K_i, \] (70)
  and as \( \epsilon \to 0 \) the contracted algebra is obtained by the substitutions
  \[ \epsilon_1 \to \epsilon_1, \quad \lambda \to \lambda, \quad \epsilon_2 \to 0. \] (71)
  As (70) indicates, the corresponding limiting case is characterized by small speeds (parameters of the inertial transformations) and small
space intervals. So this contraction can be called Speed-Space (sl)
contraction [79]. According to (71), the Speed-Space contraction de-
scribes a transition from the relative-time groups to the absolute-time
groups:

$$DS \rightarrow NH, \quad ADS \rightarrow ANH, \quad P \rightarrow G,$$

$$AP' \rightarrow G', \quad P' \rightarrow G', \quad C \rightarrow S.$$  

- Contraction with respect to the three-dimensional Euclidean group
generated by \((J_i, P_i)\), which is the motion group of the three-dimensio-
nal Euclidean space. Under this contraction

$$H \rightarrow \epsilon H, \quad K_i \rightarrow \epsilon K_i,$$

and the limit \(\epsilon \rightarrow 0\) produces the changes

$$\epsilon_1 \rightarrow \epsilon_1, \quad \lambda \rightarrow 0, \quad \epsilon_2 \rightarrow \epsilon_2.$$  

The physical meaning of this contraction is the limit when speeds are
small and time intervals are small; hence the name “Speed-Time (st)
contraction”. The speed-Time contraction leads to the absolute-space
groups with essentially no causal relations between events, and hence,
of reduced physical significance:

$$DS \rightarrow AP', \quad ADS \rightarrow P', \quad P \rightarrow C,$$

$$NH \rightarrow G', \quad ANH \rightarrow G', \quad G \rightarrow S.$$  

The absolute space groups themselves remain intact under the Speed-
Time contraction.

- Contraction with respect to the Lorentz group generated by \((J_i, K_i)\).
This is the Space-Time (lt) contraction because under this contraction

$$H \rightarrow \epsilon H, \quad P_i \rightarrow \epsilon P_i$$

and we get the limiting case of small space and time intervals. The
contracted groups are obtained by the changes

$$\epsilon_1 \rightarrow 0, \quad \lambda \rightarrow \lambda, \quad \epsilon_2 \rightarrow \epsilon_2$$

and we get transitions from the global (cosmological) groups to the
local groups:

$$DS \rightarrow P, \quad ADS \rightarrow P, \quad NH \rightarrow G, \quad ANH \rightarrow G,$$

$$AP' \rightarrow C, \quad P' \rightarrow C, \quad G' \rightarrow S.$$
• Contraction with respect to the rotation subgroup generated by $J_i$.
Under this Speed-Space-Time (slt) contraction

$$ H \rightarrow \epsilon H, \quad P_i \rightarrow \epsilon P_i, \quad K_i \rightarrow \epsilon K_i, $$

which leads in the limit $\epsilon \rightarrow 0$ to

$$ \epsilon_1 \rightarrow 0, \quad \lambda \rightarrow 0, \quad \epsilon_2 \rightarrow 0. $$

That is, all kinematical groups are contracted into the static group.

Schematically the relations between various kinematical groups are shown in Fig. 3. All these groups are limiting cases of the de Sitter or anti de Sitter groups.

As we have seen above, the natural parameter of the inertial transformations in the Poincaré group is the rapidity, which is dimensionless. One can ask the similar question about the natural dimension of speeds in other kinematical groups also. Note that the term speeds, as opposed to velocities, will be used to denote natural parameters of inertial transformations in general.

It is natural to choose group parameters for which the Lie algebra structure constants are dimensionless [88]. At that, as the Lie algebra generators have dimensions inverse to those of the corresponding group parameters, every non-zero commutation relation between $H, P_i$ and $K_i$ will induce a non-trivial relation between their dimensions $[H], [P_i]$ and $[K_i]$.

In the case of static kinematics, all commutation relations vanish. Therefore, there are no non-trivial dimensional relations, and the dimensions of time translations, space translations and speeds, which we denote respectively by $T, L$ and $S$ (so that $[H] = T^{-1}$, $[P_i] = L^{-1}$, $[K_i] = S^{-1}$), are all independent.
For the de Sitter and anti de Sitter groups, all structure constants are non-zero, which implies the following dimensional relations (note that angles, and hence $J_i$, are dimensionless)

$$T^{-1}L^{-1} = S^{-1}, \quad T^{-1}S^{-1} = L^{-1}, \quad L^{-2} = 1, \quad S^{-2} = 1, \quad L^{-1}S^{-1} = T^{-1}. $$

This is possible only if $L = T = S = 1$. Therefore the natural group parameters for de Sitter and anti de Sitter kinematics are dimensionless.

"From our 'Galilean' viewpoint, we could say that in the de Sitter universe there is a 'characteristic' length, a 'characteristic' time and a 'characteristic' speed which may be used as natural units, and then lengths, times and speeds are dimensionless" [88].

In Poincaré kinematics we have three non-trivial dimensional relations

$$T^{-1}S^{-1} = L^{-1}, \quad S^{-2} = 1, \quad L^{-1}S^{-1} = T^{-1}, $$

which imply that speeds are dimensionless and $L = T$. That is, space and time are unified in one dimensional quantity while speeds are natural to measure in terms of a characteristic speed $c$.

In Galilei kinematics $T^{-1}S^{-1} = L^{-1}$ and we get $S = LT^{-1}$, which is the usual velocity. Space and time in this case are independent dimensional quantities.

Dimensional structures of other kinematical groups are more exotic [88]. In para-Poincaré and anti para-Poincaré case $L = 1$ and $S = T$. That is, there is a characteristic length while time and speeds are unified in one dimensional quantity. In Carroll case $S = TL^{-1}$ and space and time are independent dimensional quantities. Newton–Hook and anti Newton–Hook space-times are characterized by $T = 1$ and $S = L$; that is, there is a characteristic time while length and speeds are unified. At last, in para-Galilei case space and time are independent dimensional quantities but $S = LT$.

Remarkably, all kinematical groups admit a four-dimensional space-time which can be identified with the homogeneous space of the group, namely, with its quotient by the six-dimensional subgroup generated by the rotations $J_i$ and the inertial transformations $K_i$. At that, for the kinematical groups with vanishing commutators $[K_i, H] = 0$ and $[K_i, P_j] = 0$, that is, for the para-Galilei and static groups, inertial transformations do not act on the space-time. For other groups the space-time transforms non-trivially under inertial transformations. Let us find the corresponding transformations for Newton–Hook groups, for example, in the case of $(1+1)$-dimensional space-time to avoid unnecessary technical details.

A space-time point $(x_0, x)$ is the equivalence class of the group element $e^{x_0 H} e^x P$. After the inertial transformation $e^{\psi K}$, we get a new point $(x'_0, x')$,
which, on the other hand, is the equivalence class of the group element \( e^{\psi K} e^{x_0 H} e^{x P} e^{-\psi K} \). For Newton–Hook groups \( \epsilon_2 = 0 \) and, as (48) implies, the generators \( H = G_1, P = G_2, K = G_3 \) satisfy

\[
\begin{align*}
H^2 &= \epsilon_1 E_{23}, & E_{23} H &= H, & P^2 &= 0, & K^2 &= 0, & E_{23} K &= 0, \\
E_{23} P &= 0, & H P &= 0, & H K &= 0, & P K &= 0
\end{align*}
\]

with

\[
E_{23} = \begin{pmatrix} 0 & 0 & 0 \\ 0 & 1 & 0 \\ 0 & 0 & 1 \end{pmatrix}.
\]

Using the first two of these relations, we get

\[
e^{x_0 H} = 1 + [C(x_0; \epsilon_1) - 1] E_{23} + S(x_0; \epsilon_1) H,
\]

and after applying the Baker–Campbell–Hausdorff formula (51), along with the commutation relations (61), one easily obtains

\[
e^{\psi K} e^{x_0 H} e^{-\psi K} = 1 + [C(x_0; \epsilon_1) - 1] E_{23} + S(x_0; \epsilon_1) H + \psi [C(x_0; \epsilon_1) - 1] K - \psi S(x_0; \epsilon_1) P.
\]

Because of the relations (72), the right-hand-side is the same as

\[
e^{x_0 H} e^{-\psi S(x_0; \epsilon_1)} e^{\psi [C(x_0; \epsilon_1) - 1] K}.
\]

But \( e^{\psi K} e^{x P} e^{-\psi K} = e^{x P} \), as \( P \) and \( K \) do commute, and therefore,

\[
e^{\psi K} e^{x_0 H} e^{x P} e^{-\psi K} = e^{\psi K} e^{x_0 H} e^{-\psi K} e^{\psi K} e^{x P} e^{-\psi K} = e^{x_0 H} e^{[x - \psi S(x_0; \epsilon_1)] P} e^{\psi [C(x_0; \epsilon_1) - 1] K}
\]

which has the same equivalence class as \( e^{x_0 H} e^{[x - \psi S(x_0; \epsilon_1)] P} \). Therefore, the transformation law is

\[
\begin{align*}
x'_0 &= x_0, \\
x' &= x - \psi S(x_0; \epsilon_1).
\end{align*}
\]

For \( \epsilon_1 = 0 \), we get Galilean transformations with \( x_0 = t \) and \( \psi = V \), which are the natural choice of group parameter dimensions for the Galilei kinematics.

For Newton–Hook kinematics we have a characteristic time \( \tau \) and under the natural choice of group parameters \( \epsilon_1 = 1, x_0 = t/\tau, \psi = V\tau \) the transformation law (73) takes the form

\[
\begin{align*}
x'_0 &= x_0, \\
x' &= x - V\tau \sinh \frac{t}{\tau}.
\end{align*}
\]
In contrast to the Poincaré or Galilean case, a particle at rest $x = 0$ experiences exponentially accelerated inertial motion

$$x' = -V \tau \sinh \frac{t'}{\tau}$$

if subjected to a pure inertial transformation. Therefore, Newton–Hook kinematics corresponds to expanding universe [79].

In the case of anti Newton–Hook kinematics, $\epsilon_1 = -1$ and we have oscillating universe with inertial transformations

$$x'_0 = x_0,$$
$$x' = x - V \tau \sin \frac{t}{\tau}.$$ (75)

Would Minkowski follow the logic of his staircase-wit from the "Raum und Zeit" to the end, he would reveal that the Poincaré group is mathematically less intelligible than the de Sitter or anti de Sitter groups, whose limiting cases it is. Therefore, he missed an opportunity to predict the expanding universe already in 1908 [46].

However, as the existence of the Newton–Hook groups indicate, the Newtonian ($\epsilon_2 = 0$) relativity is sufficient to predict the expanding universe and it is not well-known that the expanding universe was indeed proposed by Edgar Allan Poe in 1848, many years before Friedmann and Lemaître! Poe’s *Eureka*, from which we borrowed the preface, is an extended version of the lecture Poe gave at the Society Library of New York. The lecture was entitled “On the Cosmogony of the Universe”, and despite its quite naive and metaphysical premises this bizarre mixture of metaphysics, philosophy, poetry, and science contains several brilliant ideas central in modern-day cosmology, including a version of the Big Bang and evolving universe with inflation of the primordial atom at the start, resolution of the Olbers’ paradox (why sky is dark at night), an application of the Anthropic Cosmological Principle to explain why the universe is so immensely large, a suggestion of the multiverse with many causally disjoint universes, each with its own set of physical laws [89–91]. Although this is quite fascinating, it seems Poe was driven mainly by his poetic aesthetics in producing these ideas rather than by scientific logic [92]. As a result, these ideas, being far ahead of the time, remained obscure for contemporaries and have not played any significant role in the historical development of cosmology. Curiously, as witnessed by his biographers, Poe was Friedmann’s favorite writer. “Did Friedmann read *Eureka*? It would be not serious to push this game too far” [89].
10. Once more about mass

It seems worthwhile to return to “the virus of relativistic mass” [93] and inspect the concept of mass from the different viewpoint provided by the quantum theory. After all, the creation of quantum mechanics was another and more profound conceptual revolution in physics. Unfortunately, the corresponding dramatic change of our perspective of reality is usually ignored in teaching relativity.

The states of a quantum-mechanical system are described by vectors in a Hilbert space. At that, vectors which differ only in phase represent the same state. That is, the quantum-mechanical state is represented by a ray

\[ e^{i\alpha} |\Psi\rangle, \]

where \( \alpha \) is an arbitrary phase, rather than by a single vector \( |\Psi\rangle \). This gauge freedom has some interesting consequences for a discussion of symmetries in quantum case.

Symmetries of a quantum system are represented by unitary (or anti-unitary, if time reversal is involved) operators in the Hilbert space. Because of the gauge freedom (76), these unitary operators are also defined only up to phase factors. Let us take a closer look at these things assuming that under the symmetry transformation \( g \in G \) the space-time point \( x = (t, \vec{x}) \) transforms into \( x' = (t', \vec{x}') \). For the short-hand notation we will write \( x' = g(x) \). Let \( |x\rangle \equiv |\vec{x}, t\rangle \) be the basis vectors of the \( x \)-representation, that is, of the representation where the coordinate operator is diagonal. The symmetry \( g \) is represented in the Hilbert space by a unitary operator \( U(g) \) and, obviously, \( |x'\rangle \) and \( U(g)|x\rangle \) should represent the same state, that is, they can differ only by a phase factor:

\[ |x'\rangle = e^{i\alpha(x; g)} U(g)|x\rangle, \quad \langle x | U(g) = e^{i\alpha_1(x; g)} \langle g^{-1}(x)|. \]

(77)

Here the second identity follows from the first one when we take into account that operators \( U^+(g) = U^{-1}(g) \) and \( U(g^{-1}) \) can differ from each other only by a phase factor. For an arbitrary state vector \( |\Psi\rangle \) we have the transformation law

\[ |\Psi'\rangle = U(g)|\Psi\rangle. \]

Using (77), one has for the wave function

\[ \Psi'(\vec{x}, t) = \langle x | \Psi'\rangle = \langle x | U(g)|\Psi\rangle = e^{i\alpha_1(x; g)} \langle x'' | \psi\rangle = e^{i\alpha_1(x; g)} \Psi(\vec{x}'', t''), \]

where \( x'' = g^{-1}(x) \). Therefore,

\[ \Psi'(x) = e^{i\alpha_1(x; g)} \Psi(g^{-1}(x)). \]

(78)
Now let us compare \( \langle x|U(g_1)U(g_2)|\Psi \rangle \) and \( \langle x|U(g_1 \cdot g_2)|\Psi \rangle \) [94]. Using (77), we get
\[
\langle x|U(g_1)U(g_2)|\Psi \rangle = e^{i\alpha_1(x;g_1)} \langle x_1|U(g_2)|\Psi \rangle = e^{i\alpha_1(x;g_1)} e^{i\alpha_1(x_1;g_2)} \langle x_{12}|\Psi \rangle ,
\]
where \( x_1 = g_1^{-1}(x) \) and \( x_{12} = (g_1 \cdot g_2)^{-1}(x) \). On the other hand,
\[
\langle x|U(g_1 \cdot g_2)|\Psi \rangle = e^{i\alpha_1(x;g_1 \cdot g_2)} \langle x_{12}|\Psi \rangle .
\]
Therefore,
\[
\langle x|U(g_1)U(g_2)|\Psi \rangle = e^{i\alpha_2(x;g_1;g_2)} \langle x|U(g_1 \cdot g_2)|\Psi \rangle \tag{79}
\]
for all space-time points \( x \) and for all state vectors \( |\Psi \rangle \). Here
\[
\alpha_2(x;g_1, g_2) = \alpha_1(g_1^{-1}(x);g_2) - \alpha_1(x;g_1 \cdot g_2) + \alpha_1(x;g_1) = (\delta \alpha_1)(x;g_1,g_2). \tag{80}
\]

It will be useful to use elementary cohomology terminology [95], although we will not go into any depth in this high-brow theory. Any real function \( \alpha_n(x;g_1, g_2, \ldots, g_n) \) will be called a cochain. The action of the coboundary operator \( \delta \) on this cochain is determined as follows:
\[
(\delta \alpha_n)(x;g_1, g_2, \ldots, g_n, g_{n+1}) = \\
\alpha_n(g_1^{-1}(x);g_2, g_3, \ldots, g_n, g_{n+1}) - \alpha_n(x;g_1 \cdot g_2, g_3, \ldots, g_n, g_{n+1}) \\
+ \alpha_n(x;g_1, g_2 \cdot g_3, g_4, \ldots, g_n, g_{n+1}) - \alpha_n(x;g_1, g_2, g_3 \cdot g_4, g_5, \ldots, g_n, g_{n+1}) \\
+ \cdots + (-1)^n \alpha_n(x;g_1, g_2, \ldots, g_n \cdot g_{n+1}) + (-1)^{n+1} \alpha_n(x;g_1, g_2, \ldots, g_n). \tag{81}
\]

The coboundary operator has the following fundamental property
\[
\delta^2 = 0. \tag{82}
\]
A cochain with zero coboundary is called a cocycle. Because of (82), every coboundary \( \alpha_n = \delta \alpha_{n-1} \) is a cocycle. However, not all cocycles can be represented as coboundaries. Such cocycles will be called nontrivial.

Low dimensional cocycles play an important role in the theory of representations of the symmetry group \( G \) [95]. For example, if \( \alpha_1(x;g) \) is a cocycle, so that \( \alpha_2(x;g_1, g_2) \) vanishes, then (79) indicates that
\[
U(g_1 \cdot g_2) = U(g_1)U(g_2) \tag{83}
\]
and the unitary operators \( U(g) \) realize a representation of the group \( G \).
However, because of the gauge freedom related to the phase ambiguity, it is not mandatory that the unitary operators $U(g)$ satisfy the representation property (83). It will be sufficient to have a projective (or ray) representation.

\begin{equation}
U(g_1 \cdot g_2) = e^{-i\xi(g_1,g_2)}U(g_1)U(g_2).
\end{equation}

The so-called local exponent $\xi(g_1,g_2)$ cannot be quite arbitrary. In particular, the associativity property $(g_1 \cdot g_2) \cdot g_3 = g_1 \cdot (g_2 \cdot g_3)$ implies that

\begin{equation}
\xi(g_1,g_2) + \xi(g_1 \cdot g_2,g_3) = \xi(g_2,g_3) + \xi(g_1,g_2 \cdot g_3),
\end{equation}

and, therefore, $\xi(g_1,g_2)$ is a global (not-dependent on the space-time point $x$) cocycle:

\begin{equation}
(\delta\xi)(g_1,g_2,g_3) = \xi(g_2,g_3) - \xi(g_1 \cdot g_2,g_3) + \xi(g_1,g_2 \cdot g_3) - \xi(g_1,g_2) = 0.
\end{equation}

If the cocycle (80) does not depend on $x$, one can identify

\begin{equation}
\xi(g_1,g_2) = \alpha_2(x;g_1,g_2).
\end{equation}

At that $\alpha_2(x;g_1,g_2)$ is a trivial local cocycle, as (80) indicates, but globally it is not necessarily trivial, that is, representable as the coboundary of a global cochain.

If unitary operators $U(g)$ constitute a projective representation of the symmetry group $G$, then the correspondence $(\theta,g) \rightarrow e^{i\theta}U(g)$, where $\theta$ is a real number and $g \in G$, gives an ordinary representation of the slightly enlarged group $\tilde{G}$ consisting of all pairs $(\theta,g)$. A group structure on $\tilde{G}$ is given by the multiplication law

\begin{equation}
(\theta_1,g_1) \cdot (\theta_2,g_2) = (\theta_1 + \theta_2 + \xi(g_1,g_2), g_1 \cdot g_2).
\end{equation}

Indeed, the cocycle condition (85) ensures that the multiplication law (86) is associative. If we assume, as is usually done, that $U(e) = 1$, where $e$ is the unit element of the group $G$, then (84) indicates that $\xi(e,e) = 0$. By setting, respectively, $g_2 = g_3 = e$, $g_1 = g_2 = e$ and $g_1 = g_3 = g$, $g_2 = g^{-1}$ in (85), we get

\begin{equation}
\xi(e,g) = \xi(g,e) = 0, \quad \xi(g,g^{-1}) = \xi(g^{-1},g),
\end{equation}

for every element $g \in G$. However, it is evident then that $(0,e)$ constitutes the unit element of the extended group $\tilde{G}$ and the inverse element of $(\theta,g)$ is given by

\begin{equation}
(\theta,g)^{-1} = (-\theta - \xi(g,g^{-1}), g^{-1}).
\end{equation}
Elements of the form \((\theta, e)\) commute with all elements of \(\tilde{G}\), that is, they belong to the center of \(\tilde{G}\), and respectively \(\tilde{G}\) is called a central extension of \(G\).

The structure of the group \(G\) and the functional relation (85) greatly constrain the possible forms of admissible two-cocycles \(\xi(g_1, g_2)\) [96,97]. We will not reproduce here the relevant mathematics with somewhat involved technical details, but instead clarify the physical meaning behind this mathematical construction [94,96,98,99].

Let \(\Psi(\vec{x}, t)\) be a wave function of a free non-relativistic particle of mass \(m\). Then \(\Psi(\vec{x}, t)\) satisfies the Schrödinger equation

\[
\frac{i}{\hbar} \frac{\partial \Psi(\vec{x}, t)}{\partial t} = -\frac{1}{2\mu} \Delta \Psi(\vec{x}, t),
\]

where \(\mu = m/\hbar\).

The Galilean invariance ensures that \(\Psi'(\vec{x}, t)\) is also a solution of the same Schrödinger equation (88), because \(\Psi'\) is just the same wave function in another inertial reference frame. But according to (78)

\[
\Psi(\vec{x}, t) = \exp \left[-i\alpha_1(\vec{x}, t; g)\right] \Psi'(\vec{x}', t'),
\]

where \(x' = g(x)\), or

\[
x'_i = R_{ij}x_j - V_i t + a_i, \\
t' = t + b,
\]

\(R\) being the spatial rotation matrix: \(RR^T = R^T R = 1\). From (89) we get

\[
\frac{\partial}{\partial t} = \frac{\partial t'}{\partial t} \frac{\partial}{\partial t'} + \frac{\partial x'_i}{\partial t} \frac{\partial}{\partial x'_i} = \frac{\partial}{\partial t'} - \vec{V} \cdot \nabla'
\]

and

\[
\nabla_i = \frac{\partial x'_j}{\partial x_i} \frac{\partial}{\partial x'_j} = R_{ij} \nabla'_j, \quad \Delta = \Delta'.
\]

Therefore the Schrödinger equation (88) can be rewritten as follows (we have dropped the primes except in \(\Psi'\))

\[
\left(\frac{i}{\hbar} \frac{\partial}{\partial t} - i \vec{V} \cdot \nabla + \frac{1}{2\mu} \Delta\right) \left\{ \exp \left[-i\alpha_1(\vec{x}, t; g)\right] \Psi'(\vec{x}, t) \right\} = 0. \tag{90}
\]

But

\[
\frac{\partial}{\partial t} (e^{-i\alpha_1} \Psi') = e^{-i\alpha_1} \left[-i \frac{\partial \alpha_1}{\partial t} \Psi' + \frac{\partial \Psi'}{\partial t}\right], \\
\nabla (e^{-i\alpha_1} \Psi') = e^{-i\alpha_1} \left[-i (\nabla \alpha_1) \Psi' + \nabla \Psi'\right],
\]
and
\[
\Delta(e^{-i\alpha_1} \Psi') = e^{-i\alpha_1} \left[ -i(\Delta \alpha_1) \Psi' - (\nabla \alpha_1)^2 \Psi' - 2i(\nabla \alpha_1) \cdot (\nabla \Psi') + \Delta \Psi' \right].
\]

Therefore, (90) takes the form
\[
\Psi' \left[ \frac{\partial \alpha_1}{\partial t} - \vec{V} \cdot \nabla \alpha_1 - \frac{i}{2\mu} \Delta \alpha_1 - \frac{1}{2\mu} (\nabla \alpha_1)^2 \right]
+ \left( i \frac{\partial}{\partial t} + \frac{1}{2\mu} \Delta \right) \Psi' - i \left( \vec{V} + \frac{1}{\mu} \nabla \alpha_1 \right) \cdot \nabla \Psi' = 0.
\]

As \( \Psi' \) satisfies the Schrödinger equation (88), we are left with two more conditions on the \( \alpha_1 \) cochain:
\[
\vec{V} + \frac{1}{\mu} \nabla \alpha_1 = 0 \tag{91}
\]
and
\[
\frac{\partial \alpha_1}{\partial t} - \vec{V} \cdot \nabla \alpha_1 - \frac{i}{2\mu} \Delta \alpha_1 - \frac{1}{2\mu} (\nabla \alpha_1)^2 = 0. \tag{92}
\]

From (91), we get
\[
\alpha_1(\vec{x}, t) = -\mu \vec{V} \cdot \vec{x} + f(t)
\]
and (92) takes the form
\[
\frac{df}{dt} + \frac{1}{2} \mu V^2 = 0.
\]

Therefore,
\[
f(t) = -\frac{1}{2} \mu V^2 t + C_g,
\]
where \( C_g \) is an integration constant that can depend on \( g \), that is, on \((b, \vec{a}, \vec{V}, R)\). A convenient choice [99] is
\[
C_g = \frac{1}{2} \mu V^2 b + \mu \vec{V} \cdot \vec{a}.
\]

Then \( \alpha_1(\vec{x}, t; g) \) cochain takes the form
\[
\alpha_1(\vec{x}, t; g) = -\mu \vec{V} \cdot (\vec{x} - \vec{a}) - \frac{1}{2} \mu V^2 (t - b). \tag{93}
\]

Substituting (93) into (80), one can check that \( \xi(g_1, g_2) = \alpha_2(x; g_1, g_2) \) is indeed a global cocycle (does not depend on \( x \)):
\[
\xi(g_1, g_2) = \frac{1}{2} \mu V_1^2 b_2 - \mu \vec{V}_1 \cdot (R_1 \vec{a}_2), \tag{94}
\]
where $\vec{V}_1 \cdot (R_1 \vec{a}_2) = (R_1^{-1} \vec{V}_1) \cdot \vec{a}_2 = V_{1i} R_{1ij} a_{2j}$. Note that (94) differs from the Bargmann’s result [96] by a coboundary $C_{g_2} - C_{g_1,g_2} + C_{g_1}$ (note the overall sign difference in $\alpha_1, \alpha_2$, as well as the sign difference of the velocity term in (89) from the conventions of [96]). Cocycles which differ by a coboundary should be considered as equivalent because unitary operators $U(g)$ are determined only up to a phase. If we change representatives of the projective representation rays as follows

$$U(g) \rightarrow e^{i\phi(g)} U(g),$$

the two-cocycle $\xi(g_1,g_2)$ will be changed by a coboundary

$$\xi(g_1,g_2) \rightarrow \xi(g_1,g_2) + \phi(g_2) - \phi(g_1 \cdot g_2) + \phi(g_1).$$

In particular, different choices of the integration constant $C_g$ produce equivalent cocycles.

The Bargmann cocycle (94) is nontrivial. This can be shown as follows [96,99]. The multiplication table of the Galilei group is given by

$$(b_1, \vec{a}_1, \vec{V}_1, R_1) \cdot (b_2, \vec{a}_2, \vec{V}_2, R_2)$$

$$= (b_1 + b_2, \vec{a}_1 + R_1 \vec{a}_2 - b_2 \vec{V}_1, \vec{V}_1 + R_1 \vec{V}_2, R_1 R_2),$$

$$(b, \vec{a}, \vec{V}, R)^{-1} = (-b, -R^{-1} \vec{a} - bR^{-1} \vec{V}, -R^{-1} \vec{V}, R^{-1}), \tag{95}$$

where $R\vec{a}$ denotes the action of the rotation $R$ on the vector $\vec{a}$, that is $(R\vec{a})_i = R_{ij} a_j$. We have already used (95) while deriving (94).

It follows from this multiplication table that elements of the form $(0, \vec{a}, \vec{V}, 1)$ (space translations and boosts) form an Abelian subgroup:

$$(0, \vec{a}_1, \vec{V}_1, 1) \cdot (0, \vec{a}_2, \vec{V}_2, 1) = (0, \vec{a}_1 + \vec{a}_2, \vec{V}_1 + \vec{V}_2, 1),$$

$$(0, \vec{a}_1, \vec{V}_1, 1)^{-1} = (0, -\vec{a}, -\vec{V}, 1).$$

every trivial cocycle, having the form $\phi(g_2) - \phi(g_1 \cdot g_2) + \phi(g_1)$, is necessarily symmetric in $g_1,g_2$ on Abelian subgroups. However, the Bargmann cocycle remains asymmetric on the Abelian subgroup of space translations and boosts:

$$\xi(g_1,g_2) = -\mu \vec{V}_1 \cdot \vec{a}_2, \quad g_1 = (0, \vec{a}_1, \vec{V}_1, 1), \quad g_2 = (0, \vec{a}_2, \vec{V}_2, 1).$$

Therefore, it cannot be a trivial cocycle.

The presence of a nontrivial two-cocycle in quantum field theory usually signifies the existence of anomalous Schwinger terms in current commutators [100,101]. Remarkably, in Galilean quantum mechanics mass plays the role of the Schwinger term, as we demonstrate below [94].
Let us take $g_1 = (0, \vec{V}, 1)$ and $g_2 = (0, \vec{a}, 0, 1)$. Then $\xi(g_1, g_2) = -\mu \vec{V} \cdot \vec{a}$ and $\xi(g_2, g_1) = 0$. Therefore,

$$U(g_1)U(g_2) = e^{-i\mu \vec{V} \cdot \vec{a}}U(g_1 \cdot g_2) = e^{-i\mu \vec{V} \cdot \vec{a}}U(g_2)U(g_1).$$

(96)

In terms of the (anti-Hermitian) generators $K_i$ and $P_i$ of the Galilean group, we have

$$U(g_1) = e^{\vec{K} \cdot \vec{V}}, \quad U(g_2) = e^{\vec{P} \cdot \vec{a}},$$

and (96) takes the form

$$e^{\vec{K} \cdot \vec{V}} e^{\vec{P} \cdot \vec{a}} e^{-\vec{K} \cdot \vec{V}} e^{-\vec{P} \cdot \vec{a}} = e^{-i\mu \vec{V} \cdot \vec{a}}.$$

(97)

Expanding up to second-order terms in the infinitesimals $\vec{V}$ and $\vec{a}$, we get from (97)

$$[\vec{K} \cdot \vec{V}, \vec{P} \cdot \vec{a}] = -i\mu \vec{V} \cdot \vec{a},$$

or

$$[P_i, K_j] = i\mu \delta_{ij}.$$  

(98)

The commutator (98) is anomalous because space translations and boosts commute on the group level while their generators in the representation $U(g)$ do not.

The analogous commutator in the Poincaré group Lie algebra

$$[P_i, K_j] = \delta_{ij} H$$

(99)

is not anomalous, as the following argument shows [94]. Let $g_1$ be a pure Lorentz transformation

$$t' = \gamma \left( t - \frac{\vec{V} \cdot \vec{x}}{c^2} \right),$$

$$\vec{x}' = \vec{x} - \gamma \vec{V} t + \frac{\gamma^2}{1 + \gamma^2} \frac{\vec{V} \cdot \vec{x}}{c^2} \vec{V},$$

and $g_2$, a space translation

$$t' = t,$$

$$\vec{x}' = \vec{x} + \vec{a}.$$  

Then for infinitesimal $\vec{V}$ and $\vec{a}$ we get that $g_1 \cdot g_2 \cdot g_1^{-1} \cdot g_2^{-1}$ is a pure time translation:

$$\left( g_1 \cdot g_2 \cdot g_1^{-1} \cdot g_2^{-1} \right) \begin{pmatrix} t \\ \vec{x} \end{pmatrix} = \begin{pmatrix} t - \frac{\vec{V} \cdot \vec{a}}{c^2} \\ \vec{x} + \frac{\gamma^2}{1 + \gamma} \frac{\vec{V} \cdot \vec{a}}{c^2} \vec{V} \end{pmatrix} \approx \begin{pmatrix} t - \frac{\vec{V} \cdot \vec{a}}{c^2} \\ \vec{x} \end{pmatrix},$$

(100)
up to second-order terms. However, this is exactly what is expected from (99). For infinitesimal \( \vec{V} \), the Lorentz transformation \( g_1 \) is represented by the unitary operator

\[
U(g) = \exp \left( \frac{\vec{K} \cdot \vec{\beta}}{\beta} \right) \approx \exp \left( \frac{\vec{K} \cdot \vec{V}}{c} \right),
\]

where \( \psi \approx \tanh \psi = \beta \). Therefore,

\[
U(g_1)U(g_2)U^{-1}(g_1)U^{-1}(g_2) \approx 1 + \left[ \frac{\vec{K} \cdot \vec{V}}{c}, \vec{P} \cdot \vec{a} \right] = 1 - \frac{\vec{V} \cdot \vec{a}}{c} H.
\]

While the time translation (100) is represented by

\[
e^{\Delta x_0 H} = \exp \left( -\frac{\vec{V} \cdot \vec{a}}{c} H \right) \approx 1 - \frac{\vec{V} \cdot \vec{a}}{c} H,
\]

where \( x_0 = ct \).

The above argument indicates that the Poincaré group does not admit nontrivial cocycles. This can be confirmed by the following observation. Instead of the Schrödinger equation (88), in the relativistic case we will have the Klein–Gordon equation

\[
(\Box + \mu^2 c^2) \Psi(x) = 0,
\]

where \( \Box = \partial_{\mu} \partial^{\mu} = \Box' \) is now invariant under the action of the symmetry (Poincaré) group and, therefore, we will have, instead of (90), the following equation

\[
(\Box + \mu^2 c^2) [e^{-i\alpha_1} \Psi] = e^{-i\alpha_1} \left\{ (\Box + \mu^2 c^2) \Psi - 2i(\partial_{\mu} \alpha_1)(\partial^{\mu} \Psi) - [(\partial_{\mu} \alpha_1)(\partial^{\mu} \alpha_1) + i\Box \alpha_1] \Psi \right\} = 0.
\]

It follows then that

\[
\partial^{\mu} \alpha_1 = 0,
\]

which implies the independence of \( \alpha_1 \) from the space-time point \( x \),

\[
\alpha_1(x; g) = C_g,
\]

and therefore \( \xi(g_1, g_2) = (\delta \alpha_1)(g_1, g_2) \) is a globally trivial two-cocycle.
It is an interesting question how the nontrivial cocycle of the Galilei group arises in the process of the Poincaré group contraction. This can be clarified as follows [102]. The general Poincaré transformations have the form [103]

\[ t' = \gamma \left( t - \frac{\vec{V} \cdot \vec{R} \vec{x}}{c^2} \right) + b, \]

\[ \vec{x}' = \vec{R} \vec{x} - \gamma \vec{V} t + \frac{\gamma^2}{1 + \gamma} \frac{\vec{V} \cdot \vec{R} \vec{x}}{c^2} \vec{V} + \vec{a}. \]  

(101)

Let us take

\[ C_g = -\mu c^2 b, \]

so that \( \alpha_1(g) = C_g \) diverges as \( c^2 \to \infty \). Nevertheless, its coboundary \( \xi(g_1, g_2) = (\delta \alpha_1)(g_1, g_2) \) does have a well defined limit under the Poincaré group contraction. Indeed, it follows from (101) that

\[ b(g_1 \cdot g_2) = b_1 + \gamma_1 b_2 - \gamma_1 \frac{\vec{V}_1 \cdot \vec{R}_1 \vec{a}_2}{c^2}. \]

Therefore,

\[ (\delta \alpha_1)(g_1, g_2) = \alpha_1(g_2) - \alpha_1(g_1 \cdot g_2) + \alpha_1(g_1) = \mu c^2 (\gamma_1 - 1) b_2 - \mu \gamma_1 \vec{V}_1 \cdot \vec{R}_1 \vec{a}_2, \]

and this expression converges to the Bargmann cocycle (94) as \( c^2 \to \infty \):

\[ \mu c^2 (\gamma_1 - 1) b_2 - \mu \gamma_1 \vec{V}_1 \cdot \vec{R}_1 \vec{a}_2 \to \frac{1}{2} \mu V_1^2 b_2 - \mu \vec{V}_1 \cdot \vec{R}_1 \vec{a}_2. \]

These considerations indicate that mass plays somewhat different conceptual roles in relativistic and non-relativistic quantum theories. In Galilei invariant theory mass has a cohomological origin and appears as a Schwinger term (central extension parameter). Both in relativistic and non-relativistic quantum theories, mass plays the role of a label (together with spin) distinguishing different irreducible representations of the symmetry group and, therefore, different elementary quantum systems. Of course, this new quantum facets of mass are a far cry from what is usually assumed in the Newtonian concept of mass (a measure of inertia and a source of gravity). Nevertheless, I believe they are more profound and fundamental, while the Newtonian facets of mass are just emerging concepts valid only in a restricted class of circumstances (in non-relativistic situations and in the classical limit).

Considering mass as a Schwinger term has an important physical consequence because it leads to the mass superselection rule [96, 98]. Let us return to (97) which indicates that the identity transformation \( g_1 \cdot g_2 \cdot g_1^{-1} \cdot g_2^{-1} \) of the
Galilei group is represented by the phase factor $e^{-i\mu \vec{V} \cdot \vec{a}}$ in the Hilbert space. This is fine, except for coherent superpositions of different mass states for which (97) leads to an immediate trouble. For example, if $|\Psi_1\rangle$ and $|\Psi_2\rangle$ are two states with different masses $m_1 \neq m_2$, then

$$U(g_1)U(g_2)U^{-1}(g_1)U^{-1}(g_2)\left( |\Psi_1\rangle + |\Psi_2\rangle \right) = e^{-i\mu_1 \vec{V} \cdot \vec{a}} \left( |\Psi_1\rangle + e^{i\mu_1 - \mu_2} |\Psi_2\rangle \right)$$

which is physically different from $|\Psi_1\rangle + |\Psi_2\rangle$.

Therefore, quantum states with different masses should belong to different superselection sectors and can not be coherently superposed. In particular, there are no neutrino oscillations in Galilei invariant theory [104]. In the relativistic case we have no such prohibition of coherent superpositions of different energy states which must often also be considered as superpositions of different mass states. Experimental observation of neutrino oscillations, therefore, directly indicates that we do not live in a Galilei invariant world.

Besides, we see once again that the Galilei group is a rather singular and subtle limit of the Poincaré group, and not always correctly describes the non-relativistic limit of the Poincaré invariant quantum theory. There are real physical phenomena (like neutrino oscillations), persistent in the non-relativistic limit, which Galilei invariant quantum theory fails to describe [105]. One may expect, therefore, that the inverse road from Galilean world to the relativistic one can have some pitfalls (like relativistic mass). It is correct that historically just this adventurous road was used to reach the relativistic land. “This was so because Einstein was going from a Galilean universe to an as yet unknown one. So, he had to use Galilean concepts in his approach to a new theory” [88]. However, after a hundred years of persistent investigation of this new land maybe it will be wiser not to use the Newtonian road to the relativistic world any more and avoid pitfalls by embarking on safer pathways.

### 11. The return of æther?

A by-product of special relativity was that æther became a banished word in physics, given on worry to crackpots. However, the concept of æther has too venerable a history [106] and after a century of banishment we may ask whether it is reasonable to give up the term. After all we are still using the word ‘atom’ without attaching to it the same meaning as given initially by ancient Greeks.

It was Einstein himself who tried to resurrect æther as the general relativistic space-time in his 1920 inaugural lecture at the University of Leiden [107]:
“Recapitulating: we may say that according to the general theory of relativity space is endowed with physical qualities; in this sense, therefore, there exists æther. According to the general theory of relativity space without æther is unthinkable; for in such space there not only would be no propagation of light, but also no possibility of existence for standards of measuring rods and clocks, nor therefore any space-time intervals in the physical sense. But this æther may not be thought of as endowed with the quality characteristic of ponderable media, as consisting of parts which may be tracked through time. The idea of motion may not be applied to it.”

However, Einstein was more successful in eliminating æther than in giving it new life later. Actually all classical æther theories had got a death blow and became doomed after special relativity. But quantum mechanics added a new twist in the story. This was first realized (to my knowledge) by another great physicist Paul Dirac [108]. His argument goes as follows.

Usually it is supposed that æther is inconsistent with special relativity because it defines a preferred inertial reference frame — where the æther is at rest. In other reference frames the æther moves with some velocity and this velocity vector provides a preferred direction in space-time which should show itself in suitably designed experiments.

But in quantum mechanics the velocity of the æther is subject to uncertainty relations and usually it is not a well-defined quantity but distributed over a range of possible values according to the probabilities dictated by æther’s wave function. One can envisage a wave function (although not normalizable and hence describing an idealized state which can be approached indefinitely close but never actually realized) which makes all values of the æther’s velocity equally probable.

“We can now see that we may very well have an æther, subject to quantum mechanics and conforming to relativity, provided we are willing to consider the perfect vacuum as an idealized state, not attainable in practice. From the experimental point of view, there does not seem to be any objection to this. We must take some profound alterations in our theoretical ideas of the vacuum. It is no longer a trivial state, but needs elaborate mathematics for its description” [108].

The subsequent development of quantum field theory completely confirmed Dirac’s prophecy about the quantum vacuum. According to our modern perspective, the quantum vacuum is seething with activity of creating and destroying virtual quanta of various fields if probed locally. Therefore, it is much more like æther than empty space. Nevertheless, this new quantum æther is Lorentz invariant: it looks alike from all inertial reference frames. How is this possible? The following example of quasiclassical quantum æther of the electromagnetic field demonstrates the main points [109].
In the quasiclassical approximation the electromagnetic quantum æther can be viewed as space filled with a fluctuating electromagnetic field which by itself can be represented as a superposition of the transverse plane waves

\[ \vec{E}(\vec{x},t) = 2 \sum_{\lambda=1}^{2} \int d\vec{k} f(\omega) \cos \left( \omega t - \vec{k} \cdot \vec{x} - \theta(\vec{k},\lambda) \right) \vec{e}(\vec{k},\lambda), \]

\[ \vec{B}(\vec{x},t) = 2 \sum_{\lambda=1}^{2} \int d\vec{k} f(\omega) \cos \left( \omega t - \vec{k} \cdot \vec{x} - \theta(\vec{k},\lambda) \right) \frac{\vec{k} \times \vec{e}(\vec{k},\lambda)}{k}, \tag{102} \]

where \( \lambda \) labels different polarizations, the frequency \( \omega \) and wave vector \( \vec{k} \) are related by the relation \( \omega = ck, \quad k = |\vec{k}| \), and \( \vec{e}(\vec{k},\lambda) \) are unit polarization vectors:

\[ \vec{e}(\vec{k},\lambda) \cdot \vec{k} = 0, \quad \vec{e}(\vec{k},\lambda_1) \cdot \vec{e}(\vec{k},\lambda_2) = \delta_{\lambda_1\lambda_2}. \]

The fluctuating character of the electromagnetic field is indicated by introducing the uniformly distributed random phase \( \theta(\vec{k},\lambda) \) for which the following averages hold

\[ \langle \cos \theta(\vec{k}_1,\lambda_1) \cos \theta(\vec{k}_2,\lambda_2) \rangle = \frac{1}{2} \delta(\vec{k}_1 - \vec{k}_2) \delta_{\lambda_1\lambda_2}, \]

\[ \langle \sin \theta(\vec{k}_1,\lambda_1) \sin \theta(\vec{k}_2,\lambda_2) \rangle = \frac{1}{2} \delta(\vec{k}_1 - \vec{k}_2) \delta_{\lambda_1\lambda_2}, \]

\[ \langle \cos \theta(\vec{k}_1,\lambda_1) \sin \theta(\vec{k}_2,\lambda_2) \rangle = 0. \tag{103} \]

For quantum æther to be Lorentz invariant, the weight-function \( f(\omega) \) must have a special form which we will now find.

The electric and magnetic fields contribute equally to the energy density

\[ u = \frac{1}{8\pi} \langle E^2 + B^2 \rangle = \frac{1}{4\pi} \langle E^2 \rangle. \]

Substituting (102) and using (103) along with the decomposition

\[ \cos \left( \omega t - \vec{k} \cdot \vec{x} - \theta(\vec{k},\lambda) \right) \]

\[ = \cos \left( \omega t - \vec{k} \cdot \vec{x} \right) \cos \theta(\vec{k},\lambda) + \sin \left( \omega t - \vec{k} \cdot \vec{x} \right) \sin \theta(\vec{k},\lambda), \]

we obtain

\[ u = \frac{1}{4\pi} \int d\vec{k} f^2(\omega) = \int_0^\infty \rho(\omega) d\omega, \]
where the spectral energy-density function is
\[
\rho(\omega) = \frac{\omega^2}{c^3} f^2(\omega).
\]

Under a Lorentz transformation along the \(x\) axis, the \(\mathbf{E}\) and \(\mathbf{B}\) fields mix up. In particular,
\[
E'_x = E_x, \quad E'_y = \gamma (E_y - \beta B_z), \quad E'_z = \gamma (E_z + \beta B_y).
\] (104)

At that, transverse plane waves go into transverse plane waves with transformed frequencies and wave vectors
\[
\mathbf{E}'(\mathbf{x}', t') = \sum_{\lambda=1}^{2} \int d\mathbf{k} f(\omega) \cos \left( \omega' t' - \mathbf{k}' \cdot \mathbf{x}' - \theta(\mathbf{k}, \lambda) \right) \mathbf{e}'(\mathbf{k}, \lambda),
\] (105)

where
\[
\omega' = \gamma (\omega - V k_x), \quad k'_x = \gamma \left(k_x - \frac{V}{c^2} \omega \right), \quad k'_y = k_y, \quad k'_z = k_z.
\] (106)

Taking into account (102) and (104), we get for the primed polarization vectors
\[
\mathbf{e}'(\mathbf{k}, \lambda) = \gamma \mathbf{e} \left(1 - \beta \frac{k_x}{k}\right) + \gamma \beta \frac{\mathbf{e}_x}{k} \mathbf{k} + (1 - \gamma) \mathbf{e}_x \mathbf{i}.
\] (107)

It is not evident that the primed polarization vectors are transverse but this can be checked by an explicit calculation:
\[
\mathbf{e}'(\mathbf{k}, \lambda) \cdot \mathbf{k}' = \gamma \mathbf{e} \cdot \mathbf{k}' - \gamma \beta \mathbf{e}_x \left[k - \frac{k_y^2}{k} - \frac{k_z^2}{k}\right] - \gamma \beta \frac{k_x}{k} (\epsilon_y \epsilon_{y} + \epsilon_z \epsilon_{z})
\]
\[
= \gamma \mathbf{e} \cdot \mathbf{k}' \left(1 - \beta \frac{k_x}{k}\right) = 0.
\]

It follows from (108) that
\[
\mathbf{e}'(\mathbf{k}, \lambda) \cdot \mathbf{e}'(\mathbf{k}, \lambda) = \gamma^2 \left(1 - \beta \frac{k_x}{k}\right)^2
\]
\[
+ \epsilon_x^2 \left[\gamma^2 \beta^2 + (1 - \gamma)^2 + 2\beta \gamma \frac{k_x}{k} (1 - \gamma) + 2\gamma (1 - \gamma) \left(1 - \beta \frac{k_x}{k}\right)\right].
\]
However,
\[ \gamma^2 \beta^2 + (1 - \gamma)^2 + 2\gamma (1 - \gamma) = 0, \]
and we get
\[ \vec{\epsilon}'(\vec{k}, \lambda) \cdot \vec{\epsilon}'(\vec{k}, \lambda) = \gamma^2 \left( 1 - \beta \frac{k_x}{k} \right)^2. \]

(109)

It follows then that the energy density of the zero-point field in the primed frame is
\[ u' = \langle E'^2 \rangle = \frac{1}{4\pi} \int d\vec{k}' \gamma^2 \left( 1 - \beta \frac{k_x}{k} \right)^2 f^2(\omega). \]

(110)

But from (106)
\[ d\vec{k}' = \gamma \left( 1 - \beta \frac{k_x}{k} \right) d\vec{k}, \quad \gamma \left( 1 - \beta \frac{k_x}{k} \right) = \frac{\omega'}{\omega}, \]
and (110) can be rewritten as follows
\[ u' = \frac{1}{4\pi} \int d\vec{k}' \frac{\omega'}{\omega} f^2(\omega). \]

The Lorentz invariance demands that
\[ u' = u = \frac{1}{4\pi} \int d\vec{k} f^2(\omega) = \frac{1}{4\pi} \int d\vec{k}' f^2(\omega'), \]
where the last equation follows from the dummy character of the integration variable. Therefore,
\[ \frac{f^2(\omega')}{\omega'} = \frac{f^2(\omega)}{\omega} = \alpha, \]
where \( \alpha \) is some constant.

As we see, the Lorentz invariant quantum æther is ensured by the following spectral energy-density function
\[ \rho(\omega) = \alpha \frac{\omega^3}{c^3}. \]

The constant \( \alpha \) is not fixed by the Lorentz invariance alone but it can be determined from the elementary quantum theory which predicts \( \frac{1}{2} \hbar \omega \) zero-point energy per normal mode. The number of normal modes of the electromagnetic field with two independent transverse polarizations is \( 2 \frac{1}{(2\pi)^3} \) per momentum interval \( d\vec{k} = 4\pi \frac{\omega^2 d\omega}{c^3} \) and, therefore, the quantum theory predicts the spectral energy density of the quantum vacuum
\[ \rho(\omega) = \frac{4\pi \omega^2}{c^3} \left( \frac{2}{(2\pi)^3} \right) \frac{1}{2} \hbar \omega = \frac{\hbar \omega^3}{2\pi^2 c^3} \]
(111)
which fixes $\alpha$ at
\[
\alpha = \frac{\hbar}{2\pi^2}.
\]

The above derivation of the spectral energy density (111) of the Lorentz invariant quantum æther emphasizes the wave character of the electromagnetic field [109, 110]. Alternatively one can emphasize the photonic picture by considering the quantum electromagnetic æther as consisting of photons of all frequencies moving in all directions [111]. The result is the same: for radiation to be Lorentz invariant its intensity at every frequency must be proportional to the cube of that frequency.

It is instructive to show by explicit calculations that the Lorentz invariant zero-point æther does not lead to any drag force for bodies moving through it. We will demonstrate this in the framework of the Einstein–Hopf model. Einstein and Hopf showed in 1910 that in general there is a drag force slowing down a particle moving through stochastic electromagnetic background [112].

Let us consider a particle of mass $m$ and charge $e$ harmonically attached to another particle of mass $M \gg m$ and opposite charge $-e$. We assume following Einstein and Hopf that this dipole oscillator is immersed in a fluctuating electromagnetic field of the form (102), moves in the $x$-direction with velocity $V \ll c$, and is oriented so that oscillations are possible only along the $z$-axis.

In the rest frame of the oscillator, the equation of motion of mass $m$ looks like
\[
\ddot{z} - \Gamma \dot{z} + \omega_0^2 z = \frac{e}{m} E_z'(\vec{x}', t') \approx \frac{e}{m} E_z'(\vec{0}, t'),
\]
where we have assumed that the mass $M$ is located at the origin and the oscillation amplitude is small. In (112) the Abraham–Lorentz radiation reaction force is included with the well-known radiation damping constant [113]
\[
\Gamma = \frac{2 e^2}{3 mc^3}.
\]

The electric field $\vec{E}'(\vec{x}', t')$ in the primed frame is given by (105), and therefore, the steady-state solution of (112) should have the form
\[
z'(t') = \sum_{\lambda=1}^{2} \int d\vec{k} \left[ a(\omega') e^{i(\omega' t' - \theta(\vec{k}, \lambda))} + a^*(\omega') e^{-i(\omega' t' - \theta(\vec{k}, \lambda))} \right] \epsilon^*_z(\vec{k}, \lambda).
\]
(113)

Substituting (113) in (112), we get
\[
a(\omega') = \frac{e}{m} \frac{f(\omega)}{2} \frac{1}{\omega_0'^2 + i\Gamma \omega'^3 - \omega'^2}.
\]
(114)
In the $x$-direction the electromagnetic field exerts ponderomotive force on the dipole
\[ F'_x = e z' \frac{\partial E'_x}{\partial z'} - \frac{e}{c} \dot{z}' B'_y, \] (115)
here $\frac{\partial E'_x}{\partial z'}$ and $B'_y$ are evaluated at the origin $\vec{x}' = 0$ thanks to the assumed smallness of the oscillation amplitude.

Now it is straightforward, although rather lengthy, to calculate the average value of (115) by using
\[ \langle e^{i\theta(\vec{k}_1, \lambda_1)} e^{i\theta(\vec{k}_2, \lambda_2)} \rangle = \langle e^{-i\theta(\vec{k}_1, \lambda_1)} e^{-i\theta(\vec{k}_2, \lambda_2)} \rangle = 0, \]
\[ \langle e^{i\theta(\vec{k}_1, \lambda_1)} e^{-i\theta(\vec{k}_2, \lambda_2)} \rangle = \langle e^{-i\theta(\vec{k}_1, \lambda_1)} e^{i\theta(\vec{k}_2, \lambda_2)} \rangle = \delta_{\lambda_1, \lambda_2} \delta(\vec{k}_1 - \vec{k}_2), \]
which follow from (103). As a result, we get
\[ \langle e z' \frac{\partial E'_x}{\partial z'} \rangle = \sum_{\lambda=1}^{2} \int d\vec{k} \frac{e f(\omega)}{2} e_z'(\vec{k}, \lambda) \epsilon'_x(\vec{k}, \lambda) k'_z i [a(\omega') - a^*(\omega')] . \]
The polarization sum can be performed with the help of
\[ \sum_{\lambda=1}^{2} \epsilon'_i(\vec{k}, \lambda) \epsilon'_j(\vec{k}, \lambda) = \left( \frac{\omega'}{\omega} \right)^2 \left( \delta_{ij} - \frac{k'_i k'_j}{k'^2} \right), \]
which follows from (106), (107), and
\[ \sum_{\lambda=1}^{2} \epsilon_i(\vec{k}, \lambda) \epsilon_j(\vec{k}, \lambda) = \delta_{ij} - \frac{k_i k_j}{k^2} . \]
Besides,
\[ i [a(\omega') - a^*(\omega')] = \frac{e}{m} f(\omega) \frac{\Gamma_{\omega'3}}{(\omega_0'2 - \omega'^2)^2 + \Gamma^2 \omega'^6}, \]
and we get
\[ \langle e z' \frac{\partial E'_x}{\partial z'} \rangle = \int d\vec{k} \frac{e^2}{2m} f^2(\omega) \left( \frac{\omega'}{\omega} \right)^2 \left( -\frac{k'_i k'_j}{k'^2} \right) \frac{\Gamma_{\omega'^3}}{(\omega_0'2 - \omega'^2)^2 + \Gamma^2 \omega'^6} . \]
However,
\[ f^2(\omega) = \frac{e^2}{\omega^2} \rho(\omega), \quad \frac{e}{m} = \frac{3}{2} \frac{\Gamma c^3}{\omega}, \quad \frac{\omega'}{\omega} d\vec{k} = d\vec{k}', \]
and, therefore,
\[
\langle e_z' \frac{\partial E'_x}{\partial z'} \rangle = \int d\vec{k}' \frac{3}{4} \frac{e^3}{\omega^3} \rho(\omega) \left( -\frac{k'_x k'_z}{k'^2} \right) \frac{\Gamma^2 e^3 \omega'^3}{(\omega'^2 - \omega'^2)^2 + \Gamma^2 \omega'^6}.
\]

In order to change \(\rho(\omega)/\omega^3\) over to a function of \(\omega'\), we expand up to the first order in \(\beta\) [109]:
\[
\rho(\omega) \approx \rho(\omega') + \frac{\partial \rho(\omega')}{\partial \omega'} (\omega - \omega' \approx \rho(\omega') + \beta \frac{k'_x}{k'} \omega' \frac{\partial \rho(\omega')}{\partial \omega'},
\]
\[
\frac{1}{\omega^3} \approx \frac{1}{\omega'^3} - \frac{3}{\omega'^4} (\omega - \omega' \approx \frac{1}{\omega'^3} \left(1 - 3 \beta \frac{k'_x}{k'} \right).
\]

Hence, up to the first order in \(\beta\),
\[
\langle e_z' \frac{\partial E'_x}{\partial z'} \rangle = \int d\vec{k}' \frac{3}{4} \frac{e^3}{\omega'^2} \frac{\Gamma^2 e^3 \omega'^3}{(\omega'^2 - \omega'^2)^2 + \Gamma^2 \omega'^6}
\]
\[
\times \left[ \rho(\omega') + \beta \frac{k'_x}{k'} \omega' \frac{\partial \rho(\omega')}{\partial \omega'} \right] \left(3 \beta \frac{k'_x k'^2}{k'^3} - \frac{k'^2 k'_x}{k'^2}\right).
\]

Analogously, we get
\[
\langle -\frac{e}{c} \cdot \epsilon' B'_y \rangle = -\sum_{\lambda=1}^2 \int d\vec{k} \frac{e f(\omega)}{2c} \omega' \epsilon'_z \frac{(\vec{k}' \times \epsilon')_y}{k'} \epsilon'_x \left[a(\omega') - a^*(\omega')\right],
\]

because \(B'_y = \gamma (B_y + \beta E_z)\), and one can check the identity
\[
\gamma \beta \epsilon_z + \gamma \frac{(\vec{k} \times \epsilon)_y}{k} = \frac{(\vec{k}' \times \epsilon')_y}{k'}.\]

After performing the polarization sum, we are left with the expression
\[
\langle -\frac{e}{c} \cdot \epsilon' B'_y \rangle = \int d\vec{k} \frac{3}{4} \frac{e^3}{\omega^3} \rho(\omega) k'_x \frac{\Gamma^2 e^3 \omega'^3}{(\omega'^2 - \omega'^2)^2 + \Gamma^2 \omega'^6},
\]
or up to the first order in \(\beta\):
\[
\langle -\frac{e}{c} \cdot \epsilon' B'_y \rangle = \int d\vec{k}' \frac{3}{4} \frac{e^3}{\omega'^2} \frac{\Gamma^2 e^3 \omega'^3}{(\omega'^2 - \omega'^2)^2 + \Gamma^2 \omega'^6}
\]
\[
\times \left[ \rho(\omega') + \beta \frac{k'_x}{k'} \omega' \frac{\partial \rho(\omega')}{\partial \omega'} \right] \left(k'_x - 3 \beta \frac{k'^2}{k'}\right).
\]
Angular integrations, assumed in the decomposition $d\vec{k}' = \frac{\omega'^2 d\omega'}{c^3} d\Omega'$, are straightforward if the change to the polar coordinates is made and give

\[
\int \left( k_x' - 3\beta \frac{k_x'^2}{k'} + 3\beta \frac{k_x'^2 k_z'^2}{k'^3} - \frac{k_x' k_z'^2}{k'^2} \right) d\Omega' = -\frac{16\pi}{5} \beta k',
\]

\[
\int \left( k_x' - 3\beta \frac{k_x'^2}{k'} + 3\beta \frac{k_x'^2 k_z'^2}{k'^3} - \frac{k_x' k_z'^2}{k'^2} \right) \frac{k'_x}{k'} d\Omega' = \frac{16\pi}{15} k'.
\]

Therefore,

\[
\langle F'_x \rangle = -\frac{12\pi}{5} Vc \int_0^\infty \frac{\Gamma^2 \omega'^4}{(\omega_0'^2 - \omega'^2)^2 + \Gamma^2 \omega'^6} \left[ \rho(\omega') - \frac{1}{3} \omega' \frac{\partial \rho(\omega')}{\partial \omega'} \right] d\omega'.
\]  

(116)

We assume that $\Gamma \omega_0'^4 \ll 1$. But in the limit $\Gamma \omega_0' \to 0$ one has

\[
\frac{\Gamma \omega'^4}{(\omega_0'^2 - \omega'^2)^2 + \Gamma^2 \omega'^6} \to \frac{\Gamma \omega_0'^4}{(\omega_0'^2 - \omega'^2)^2 + \Gamma^2 \omega_0'^6}
\]

and

\[
\frac{\Gamma \omega_0'^3}{(\omega_0'^2 - \omega'^2)^2 + \Gamma^2 \omega_0'^6} \to \pi \delta (\omega'^2 - \omega_0'^2)
\]

\[
= \frac{\pi}{2\omega_0'} \left[ \delta(\omega' - \omega_0') + \delta(\omega' + \omega_0') \right].
\]

This makes the integration in (116) trivial, and if we drop primes in the result, because we are interested only in the lowest order terms, we finally get the expression [109, 112]

\[
\langle F_x \rangle = -\frac{6\pi^2}{5} Vc \Gamma \left[ \rho(\omega_0) - \frac{1}{3} \omega_0 \frac{\partial \rho(\omega_0)}{\partial \omega_0} \right],
\]  

(117)

for the Einstein–Hopf drag force in the laboratory frame. As expected, this drag force disappears for the Lorentz invariant æther due to the cubic dependence of its spectral energy density on the frequency (111). Alternatively one can consider the above calculations as yet another derivation of the spectral energy density which ensures the Lorentz invariance.

Interestingly, the Einstein–Hopf drag force does not vanish for the cosmic microwave background radiation with its black body spectrum

\[
\rho(\omega, T) = \frac{\hbar \omega^3}{\pi^2 c^3} \frac{1}{e^{\hbar \omega/kT} - 1}.
\]
Therefore, one has a curious situation that formally the Aristotelian view of motion is realized instead of Newtonian one: a body begins to slow down with respect to a reference frame linked to the isotropic microwave background radiation if no force is applied to it [114]. This example shows clearly that it is futile to ascribe absolute truth to physical laws. Any valid physical theory is just some idealization of reality based on concepts which are completely sound and useful within the realm of their applicability but which can go completely astray if pushed outside of this domain.

Quantum theory inevitably leads to fluctuating quantum vacuum which can be considered as a Lorentz invariant æther. The Lorentz invariance ensures that this æther does not define a preferred inertial frame and therefore the Principle of Relativity is not violated. Nevertheless, this quantum æther is completely real as it leads to experimentally observable physical effects such as mass and charge renormalizations, Casimir effect, Lamb shift, Van der Waals forces and fundamental linewidth of a laser [115]. A coherent description of these effects is provided by quantum electrodynamics. However, many aspects can qualitatively be understood even at the classical level if the existence of Lorentz invariant fluctuating æther is postulated. Probably it will come as a surprise for many physicists that the resulting classical theory, the so-called stochastic electrodynamics pioneered by Marshall [110,116] and Boyer [109,117], provides a classical foundation for key quantum concepts [118]. Of course stochastic electrodynamics can not be considered as a full-fledged substitute for the quantum theory, but it is truly remarkable that the introduction of Lorentz invariant fluctuating æther is sufficient to grasp the essence of many concepts thought to be completely quantum. The stochastic electrodynamics offers a new and useful viewpoint narrowing a gap between quantum weirdness and our classical intuition.

The electromagnetic field in the quantum æther is bound to fluctuate around the zero mean value in order to preserve the Lorentz invariance. There is no such constraint for scalar fields (elementary or composite) and they can develop non-zero vacuum expectation values. The corresponding vacuum condensates represent another example of quantum mechanical æther. These æther states play an important role in modern elementary particle theory as they lead to the phenomena of spontaneous symmetry breaking and generation of mass via the Higgs mechanism [119].

In 1993, the UK Science Minister, William Waldegrave, challenged physicists to produce a one-page answer to the question ‘What is the Higgs boson, and why do we want to find it?’ David Miller from University College London won a bottle of champagne for a very picturesque description of the Higgs mechanism [120] reproduced with slight modifications in [121].
Imagine a conference hall crowded by physicists. The physicists represent a non-trivial medium (æther) permeating the space. A gorgeous woman enters the hall and tries to find her way through the crowd. A cluster of her admirers is immediately formed around her slowing her scientific progress. At the same time the cluster gives her more momentum for the same speed of movement across the room and once moving she is harder to stop. Therefore this ephemeral creature acquires much greater effective mass. This is the Higgs mechanism. When the woman leaves the hall, a gossip about her propagates in the opposite direction bringing an excitement in the crowd. This excitement of the medium also propagates in the form of a scoundrels cluster. This is the Higgs boson which will be hunted at LHC [121].

According to Matvei Bronstein, each epoch in the history of physics has its own specific æther [122]. “The æther of the 21-st century is the quantum vacuum. The quantum æther is a new form of matter. This substance has very peculiar properties strikingly different from the other forms of matter (solids, liquids, gases, plasmas, Bose condensates, radiation, etc.) and from all the old æthers” [122].

However, there is a serious problem with the Lorentz invariant quasiclassical æther with the spectral energy density (111) because the integral

\[ u = \int_{0}^{\infty} \rho(\omega) \, d\omega \tag{118} \]

severely diverges. This problem is not cured by the full machinery of quantum field theory. It just hides this and some other infinities in several phenomenological constants, if the theory is renormalizable (like quantum electrodynamics).

The problem is possibly caused by our ignorance of the true physics at very small scales (or, what is the same, at very high energies). A natural ultraviolet cut-off in (118) is provided by the Planck frequency

\[ \omega_p = \frac{m_p c^2}{\hbar}, \quad m_p = \sqrt{\frac{\hbar c}{G}}, \]

\( G \) being the Newtonian gravitational constant. Ignoring the factors of \( 2\pi \) and the like, the particle Compton wavelength \( \hbar c/E_p \) becomes equal to its gravitational radius \( G/E_p c^4 \) at the Planck energy \( E_p = \hbar \omega_p \approx 1.22 \times 10^{19} \text{ GeV} \).

Hence the particle becomes trapped in its own gravitational field and cardinal alteration of our notions of space-time is expected [123,124].

If we accept the Planck frequency as the ultraviolet cut-off in (118), the vacuum energy density becomes (ignoring again some numerical factors
which are not relevant for the following)
\[ u \sim \frac{E_p^4}{\hbar^3 c^3} = \frac{e^7}{\hbar G^2}. \tag{119} \]

And now we have a big problem: (119) implies the cosmological constant which is fantastically too large (123 orders of magnitude!) compared to the experimental value inferred from the cosmological observations [122, 125].

Of course one can remember Dirac’s negative energy sea of the fermionic quantum fields and try to compensate (119). We could even succeed in this, thanks to supersymmetry. However, the supersymmetry, if it exists at all, is badly broken in our low energy world. Therefore it can reduce (119) somewhat but not by 123 orders of magnitude. Anyway it is much harder to naturally explain so small non-zero vacuum energy density than to make it exactly zero. This is the notorious cosmological constant problem. Among suggested solutions of this problem [126] the most interesting, in my opinion, is the one based on the analogy with condensed matter physics [125].

“Is it really surprising in our century that semiconductors and cosmology have something in common? Not at all, the gap between these two subjects practically disappeared. The same quantum field theory describes our Universe with its quantum vacuum and the many-body system of electrons in semiconductors. The difference between the two systems is quantitative, rather than qualitative: the Universe is bigger in size, and has an extremely low temperature when compared to the corresponding characteristic energy scale” [122].

The temperature that corresponds to the Planck energy is indeed very high in ordinary units \( T_p = E_p/k_B \approx 1.4 \times 10^{32} \) K. The natural question is then why all degrees of freedom are not frozen out at the temperature \( T = 300 \) K at which we live [122]. Be the Universe like ordinary semiconductors or insulators, the expected equilibrium density of excitations (elementary particles) at our living temperature \( T \) would be suppressed by the extremely large factor \( e^{T_p/T} = e^{10^{30}} \). That we survive such a freezer as our Universe indicates that the Universe is more like metals with a Fermi surface and gapless electron spectra, or, to be more precise, is like special condensed matter systems with topologically-protected Fermi points [122]. There are only a very few such systems like superfluid phases of liquid \(^3\)He and semiconductors of a special type which can be used to model cosmological phenomena [127, 128]. At that, in the condensed matter system the role of the Planck scale is played by the atomic scale. At energies much lower than the atomic scale, all condensed matter systems of Fermi point universality class exhibit a universal generic behavior and such ingredients of the Standard Model as relativistic Weyl fermions, gauge fields and effective gravity naturally emerge [127, 128]. Amusingly, if the Fermi point
topology is the main reason why the elementary particles are not frozen out at temperatures $T \ll T_p$, then we owe our own existence to the hairy ball theorem of algebraic topology which says that one cannot comb the hair on a ball smooth, because Fermi point is the hedgehog in momentum space and its stability is ensured by just that theorem [122].

Therefore, if the condensed matter analogy is really telltale, we are left with the exciting possibility that physics probably does not end even at the Planck scale and we have every reason to restore the word 'æther' in the physics vocabulary. At the present moment we can not even guess what the physics of this trans-Planckian æther is like because our familiar physical laws, as emergent low energy phenomena, do not depend much on the fine details of the trans-Planckian world, being determined only by the universality class, which the whole system belongs to. “The smaller is our energy, the thinner is our memory on the underlying high-energy trans-Planckian world of the quantum æther where we all originated from. However, earlier or later we shall try to refresh our memory and concentrate our efforts on the investigation of this form of matter” [122].

12. Concluding remarks

Although the idea of a four-vector can be traced down to Poincaré [129], it was Minkowski who gave the formulation of special relativity as a four-dimensional non-Euclidean geometry of space-time. Very few physicists were well trained in pure mathematics in general, and in non-Euclidean geometry in particular, at that time. Minkowski himself was partly motivated in his refinement of Einstein’s work, which led him to now standard four-dimensional formalism, by his doubts in Einstein’s skills in mathematics [129]. When Minkowski presented his mathematical elaboration of special relativity at a session of the Gottingen Mathematical Society, he praised the greatness of Einstein’s scientific achievement, but added that “the mathematical education of the young physicist was not very solid which I am in a good position to evaluate since he obtained it from me in Zurich some time ago” [130].

In light of Minkowski’s assessment of Einstein’s skills in mathematics, it is fair to say that “in the broader context of education in German-speaking Europe at the end of the nineteenth century Einstein received excellent preparation for his future career” [131]. Neither was he dull in mathematics as a secondary school pupil. For example, at the final examination at Aarau trade-school Einstein gave a very original solution of the suggested geometrical problem by using a general identity for the three angles of a triangle (which I am not sure is known to every physicist)

$$\sin^2 \frac{\alpha}{2} + \sin^2 \frac{\beta}{2} + \sin^2 \frac{\gamma}{2} + 2 \sin \frac{\alpha}{2} \sin \frac{\beta}{2} \sin \frac{\gamma}{2} = 1$$
and by solving a cubic equation which he got from this equality by substitutions. “Although it depended on instant recall of complicated mathematical formulas, Einstein’s solution was the very opposite of one based on brute-force calculations. He was careful to arrive at numerical values only after having made general observations on, among other things, the rationality of the roots of the cubic equation and on the geometrical requirements that a solution would have to satisfy” [131].

Anyway, after Minkowski’s sudden death from appendicitis at age 44, shortly after his seminal Cologne lecture, neither Einstein nor any other physicist at that time was in a position to duly appreciate the non-Euclidean readings of special relativity. Although the space-time formalism, energetically promoted by Sommerfeld, quickly became a standard tool in special relativity, its non-Euclidean facets remained virtually unnoticed, with just one exception.

Inspired by Sommerfeld’s interpretation of the relativistic velocity addition as a trigonometry on an imaginary sphere, the Croatian mathematician Vladimir Varičak established that relativistic velocity space possessed a natural hyperbolic (Lobachevsky) geometry [132, 133].

One can check the hyperbolic character of metric in the relativistic velocity space as follows. The natural distance between two (dimensionless) velocities \( \vec{\beta}_1 \) and \( \vec{\beta}_2 \) in the velocity space is the relative velocity [134, 135]

\[
\beta_{\text{rel}} = \sqrt{\left( \vec{\beta}_1 - \vec{\beta}_2 \right)^2 - \left( \vec{\beta}_1 \times \vec{\beta}_2 \right)^2} / \left( 1 - \vec{\beta}_1 \cdot \vec{\beta}_2 \right).
\]

Taking \( \vec{\beta}_1 = \vec{\beta} \) and \( \vec{\beta}_2 = \vec{\beta} + d\vec{\beta} \), one gets for the line element in the velocity space [134, 135]

\[
ds^2 = \frac{(d\vec{\beta})^2 - (\vec{\beta} \times d\vec{\beta})^2}{(1 - \vec{\beta}^2)^2}.
\]  

(120)

To make contact with previous formulas, we will assume two-dimensional velocity space and change to the new coordinates [136]

\[
\beta_x = \frac{2x}{1 + x^2 + y^2}, \quad \beta_y = \frac{2y}{1 + x^2 + y^2}.
\]

Then

\[
\vec{\beta} = \frac{2x}{1 + x^2 + y^2} \hat{i} + \frac{2y}{1 + x^2 + y^2} \hat{j} = f(x, y) \hat{i} + g(x, y) \hat{j},
\]

\[
d\vec{\beta} = \left( \frac{\partial f}{\partial x} dx + \frac{\partial f}{\partial y} dy \right) \hat{i} + \left( \frac{\partial g}{\partial x} dx + \frac{\partial g}{\partial y} dy \right) \hat{j},
\]
and substituting this into (120), we get after some algebra

\[ ds^2 = \frac{4(dx^2 + dy^2)}{(1 - x^2 - y^2)^2} . \]

However, this is nothing but the line element of the hyperbolic geometry from (43), with \( \epsilon_1 = 1, \epsilon_2 = -1 \).

The non-Euclidean style in relativity was pursued mainly by mathematicians and led to very limited physical insights, if any, in the relativity mainstream [132]. In a sense, non-Euclidean readings of relativity were ahead of time as illustrated by the unrecognized discovery of the Thomas precession in 1913 by the famous French mathematician Émile Borel, a former doctoral student of Poincaré.

In Borel’s non-Euclidean explanation, the Thomas precession, which is usually regarded as an obscure relativistic effect, gets a very transparent meaning. If a vector is transported parallel to itself along a closed path on the surface of a sphere, then its orientation undergoes a change proportional to the enclosed area. The Lobachevsky space is a space of constant negative curvature and, as Borel remarked, the similar phenomenon should take place: if a velocity vector circumscribes a closed path under parallel transport in the kinematical space, it will undergo precession with magnitude proportional to the enclosed area. Borel “was careful to point out that the effect is a direct consequence of the structure of the Lorentz transformations” [132]. This remarkable discovery of Borel, however, was of limited historical value because its physical significance was not recognized (I doubt it could have been done before the advent of quantum mechanics).

The hyperbolic geometry of the relativistic velocity space is an interesting but only a particular aspect of special relativity. A more important and really non-Euclidean reading of Minkowski identifies the Minkowski geometry itself as a kind of non-Euclidean geometry, as just one representative of the whole family of non-Euclidean Cayley–Klein geometries. Placing special relativity in such a broader context, we see instantly (see Fig. 3) that it is not a fundamental theory but a limiting case of more general theory. A slight generalization of the Relativity Principle (in the sense of Bacry and Lévy-Leblond [79]) leads to eleven different relativity theories which all are various limiting cases of the two really fundamental homogeneous space-times — de Sitter and anti de Sitter spaces.

It is not surprising, therefore, that astrophysical data indicate the non-zero cosmological constant, which means that the correct asymptotic (vacuum) space-time in our Universe is the de Sitter space-time, not Minkowski. What is surprising is the incredible smallness of the cosmological constant which makes the special relativity valid for all practical purposes.
As yet we do not know the resolution of the cosmological constant problem. Maybe this enigma leads to a trans-Planckian æther with yet unknown physics, as the condensed matter analogy [122] indicates. Anyway Einstein was lucky enough to be born in the Universe with nearly vanishing cosmological constant and nowadays we formulate all our theories of fundamental interactions in the Minkowskian background space-time.

The case of gravity needs some comment because it is usually assumed that special relativity is no longer correct in the presence of gravity, and should be replaced by general relativity. However, the interpretation of gravity as a pseudo-Riemannian metric of curved space-time is not the only possible interpretation, nor is it always the best one, because this interpretation sets gravity apart from other interactions — “...too great an emphasis on geometry can only obscure the deep connections between gravitation and the rest of physics” [137].

Is it possible to develop a theory of gravity as a quantum theory of massless spin-two field in the flat Minkowski space-time by analogy with other interactions? It seems Robert Kraichnan, the only post-doctoral student that Einstein ever had, was the first who initiated the study of this question. “He recalls that, though he received some encouragement from Bryce DeWitt, very few of his colleagues supported his efforts. This certainly included Einstein himself, who was appalled by an approach to gravitation that rejected Einstein’s own hard-won geometrical insights” [138]. Maybe just because of this lack of support, Kraichnan left the Institute of Advanced Study (Princeton) and the field of gravitation in 1950, to establish himself in years that followed as one of the world’s leading turbulence theorists.

Kraichnan published his work [139] only in 1955, after Gupta’s paper [140] on the similar subject appeared. Since then the approach was developed in a number of publications by various authors. It turned out that the flat Minkowski metric is actually unobservable and the geometrical interpretation of gravity as a curved and dynamical effective metric arises at the end all the same. “The fact is that a spin-two field has this geometrical interpretation: this is not something readily explainable — it is just marvelous. The geometrical interpretation is not really necessary or essential to physics. It might be that the whole coincidence might be understood as representing some kind of invariance” [141].

Although the field theory approach to gravity has an obvious pedagogical advantages [141,142], especially for high-energy physics students, we can not expect from it the same success as from quantum electrodynamics, because the theory is non-renormalizable. Besides the condensed matter analogy shows [122,143] that gravity might be really somewhat different from other interactions, and in a sense more classical, because in such interpretation gravity is a kind of elasticity of quantum vacuum (trans-Planckian æther) — the idea that dates back to Sakharov [144].
An interesting hint that gravity might be an emergent macroscopic phenomenon with some underlying microscopic quantum theory is the fact that general relativity allows the existence of space-time horizons with well-defined notions of temperature and entropy which leads to an intriguing analogy between the gravitational dynamics of the horizons and thermodynamics. It is even possible to obtain the Einstein equations from the proportionality of entropy to the horizon area together with the fundamental thermodynamical relation between heat, entropy, and temperature. “Viewed in this way, the Einstein equation is an equation of state” [145]. This perspective has important consequences for quantization of gravity: “it may be no more appropriate to canonically quantize the Einstein equation than it would be to quantize the wave equation for sound in air” [145].

All these considerations indicate that Einstein was essentially right and the Riemannian trend in geometry, which at first sight appears entirely different from Klein’s Erlangen program, is crucial for describing gravity, although this geometry might be a kind of low energy macroscopic emergent illusion and not the fundamental property of trans-Planckian æther.

Anyway, it seems reasonable to explore also another way to unify gravity with other interactions: the geometrization of these interactions. From the beginning we have to overcome a crucial obstacle on this way: for interactions other than gravity the Kleinian geometry seems to be of much more importance than the Riemannian geometry of curved space-time. “There is hardly any doubt that for physics special relativity theory is of much greater consequence than the general theory. The reverse situation prevails with respect to mathematics: there special relativity theory had comparatively little, general relativity theory very considerable, influence, above all upon the development of a general scheme for differential geometry” [146].

Just this development of mathematics, spurred by general relativity, led finally to the resolution of the dilemma we are facing in our attempts to find a common geometrical foundation for all interactions. This dilemma was formulated by Cartan as follows [55, 147]:

“The principle of general relativity brought into physics and philosophy the antagonism between the two leading principles of geometry due to Riemann and Klein respectively. The space-time manifold of classical mechanics and of the principle of special relativity is of the Klein type, and the one associated with the principle of general relativity is Riemannian. The fact that almost all phenomena studied by science for many centuries could be equally well explained from either viewpoint was very significant and persistently called for a synthesis that would unify the two antagonistic principles”.
The crucial observation which enables the synthesis is that "most properties of Riemannian geometry derive from its Levi–Civita parallelism" [53]. Let us imagine a surface in Euclidean space. A tangent vector at some point of the surface can be transferred to the nearby surface point by parallel displacement in the ambient Euclidean space, but at that, in general, it will cease to be tangent to the surface. However, we can split the new vector into its tangential and (infinitesimal) normal components and throw away the latter. This process defines the Levi–Civita connection which may be viewed as a rule for parallel transport of tangent vectors on the surface [146].

There is a more picturesque way to explain how Levi–Civita parallel transport reveals the intrinsic geometry of the surface [148]. Imagine a hamster closed inside a "hamster ball". When hamster moves inside the ball, the latter rolls on a lumpy surface. We can say that the hamster studies the intrinsic geometry of the lumpy surface by rolling more simple model surface (the sphere) on it. Rolling the ball without slipping or twisting along two different paths connecting given points on the surface will give, in general, results differing by some SO(3) rotation, an element of the principal group of the model geometry, which encodes information about the intrinsic geometry of the surface. An infinitesimal SO(3) rotation of the hamster ball, as it begins to move along some path, breaks up into a part which describes the SO(2) rotation of the sphere around the axis through the point of tangency and into a part which describes an infinitesimal translation of the point of tangency. The SO(3) connection (the analog of the Levi-Civita connection for this example) interrelates these two parts and thus defines a method of rolling the tangent sphere along the surface [148].

The far reaching generalization due to Élie Cartan is now obvious: we can take any homogeneous Klein space as a model geometry and try to roll it on the space the geometry of which we would like to study. Information about the geometry of the space under study will then be encoded in the Cartan connection which gives a method of the rolling. Figure 4 (from [148]) shows how the Cartan geometry unifies both Kleinian and Riemannian trends in geometry.

There is no conceptual difference between gravity and other physical interactions as far as the geometry is involved because the Cartan geometry provides the fully satisfying way in which gauge theories can be truly regarded as geometry [149]. Therefore, we see that special relativity and principles it rests upon are really fundamental, at the base of modern physics with its gauge theories and curved space-times. I believe its teaching should include all the beauty and richness behind it which was revealed by modern physics, and avoid historical artifacts like the second postulate and relativistic mass. I hope this article shows that such a presentation requires some quite elementary
knowledge of basic facts of modern mathematics and quantum theory. “The presentation of scientific notions as they unfolded historically is not the only one, nor even the best one. Alternative arguments and novel derivations should be pursued and developed, not necessarily to replace, but at least to supplement the standard ones” [22].

The game of abstraction outlined above is not over. For example, one hardly questions our *a priori* assumption that physical quantities are real-valued. But why? Maybe the root of this belief lies in the pre-digital age assumption that physical quantities are defined operationally in terms of measurements with classical rulers and pointers that exist in the classical continuum physical space [150]. But quantum mechanics teaches us that the Schrödinger’s cat is most likely object-oriented [151]. That suggests somewhat Platonic view of reality that physics is a concrete realization, in the realm of some Topos, of abstract logical relations among elements of reality. “Topos” is a concept proposed by Alexander Grothendieck, one of the most brilliant mathematicians of the twentieth century, as the ultimate generalization of the concept of space. Our notion of a smooth space-time manifold, upon which the real-valuefulness of the physical quantities rests, is, most likely, an emergent low energy concept and we do not know what kind of abstractions will be required when we begin “to refresh our memory” about primordial trans-Planckian æther [150, 152].

Grothendieck abruptly ended his academic career at the age of 42 and “withdrew more and more into his own tent”. If his last time visitors can be believed, “he is obsessed with the Devil, whom he sees at work everywhere in the world, destroying the divine harmony and replacing 300,000 km/sec by 299,887 km/sec as the speed of light!” [153]. I am afraid a rather radical break of the traditional teaching tradition of special relativity is required to restore its full glory and divine harmony.
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