

Random triangles and polygons in the plane

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We consider the problem of finding the probability that a random triangle is obtuse, which was first raised by Lewis Carroll. Our investigation leads us to a natural correspondence between plane polygons and the Grassmann manifold of 2-planes in real n -space proposed by Allen Knutson and Jean-Claude Hausmann. This correspondence defines a natural probability measure on plane polygons. In these terms, we answer Carroll's question. We then explore the Grassmannian geometry of planar quadrilaterals, providing an answer to Sylvester's four-point problem, and describing explicitly the moduli space of unordered quadrilaterals. All of this provides a concrete introduction to a family of metrics used in shape classification and computer vision.

The issue of choosing a “random triangle” is indeed problematic. I believe the difficulty is explained in large measure by the fact that there seems to be no natural group of transitive transformations acting on the set of triangles.

–Stephen Portnoy

A Lewis Carroll pillow problem: Probability of an obtuse triangle
Statistical Science, 1994

In 1895, the mathematician Charles L. Dodgson, better known by his pseudonym Lewis Carroll, published a book of 72 mathematical puzzles called “pillow problems”, which he claimed to have solved while lying in bed. The pillow problems mostly concern discrete probability, but there is a single problem in continuous probability in the collection:

Three points are taken at random on an infinite plane. Find the chance of their being the vertices of an obtuse-angled triangle.

This is a very appealing problem and a number of authors have tackled it in the years since. After a moment's thought, it is clear that the main issue here is that the problem is ill-posed— since there is no translation-invariant probability distribution on the infinite plane, the problem must really refer to a natural probability distribution on the space of triangles. But what probability distribution on triangle space is the right one? Portnoy [23] presented several different solutions to the problem involving distributions on triangle space invariant under various groups of transformations; Edelman and Strang [9] connect the problem to random matrix theory and shape statistics; Guy [10] got the answer $3/4$ for a variety of measures, and the legendary statistician David Kendall got exact answers when the vertices of the triangle were chosen at random in a convex body [16]. Interestingly,

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Carroll himself gave a solution, but his method gives two different answers under the assumptions that side AB is the longest or second-longest side of the triangle!

In fact, this problem has an even earlier history. In 1861, the actuary and editor of the *Lady's and Gentleman's Diary* W. S. B. Woolhouse posed the same problem for triangles in space [29]. Readers, including Stephen Watson [27] came up with the answers rediscovered by Hall [11] for triangles whose vertices are uniformly chosen in a disk or ball. In 1865 Woolhouse posed a related problem [30]:

Three lines being drawn at random on a plane, determine the probability that they will form an acute triangle.

With its focus on the edges of the triangle rather than the vertices, this version of the problem is more closely related to the approach we present in this paper, which is based on a highly symmetric representation of triangle space as a Grassmann manifold [3]. We will see that the Grassmannian picture of triangle space really does have a very natural group of transformations, that the pillow problem has a natural answer in our terms,¹ and that this entire story generalizes to the study of polygons with an arbitrary number of edges.

1. TWO PATHS TO A CONSTRUCTION

We start by fixing notation. As is usual in triangle geometry; we let A, B, C refer to the vertices of a triangle, a, b , and c denote the lengths of the corresponding (opposite) sides, and use α, β, γ for the corresponding angles.

We now construct a measure on triangle space. Since the geometry of a triangle is determined by a, b, c , it is immediately natural to want to assign a measure to the positive orthant of triples $a \geq 0, b \geq 0, c \geq 0$ in \mathbb{R}^3 . But this space is not compact, so one is led to fix a scale for the triangle, by assuming² that the perimeter $a + b + c = 2$.

Here we diverge from the beaten path. Since a, b , and c must obey the triangle inequalities $a + b \geq c, b + c \geq a$, and $c + a \geq b$, the space of triangles is actually a triangular cone inside the positive orthant. This motivates us to write things in terms of the new variables

$$s_a = \frac{-a + b + c}{2}, \quad s_b = \frac{a - b + c}{2}, \quad s_c = \frac{a + b - c}{2}.$$

These variables have a long history in triangle geometry. Perhaps most naturally if we construct three mutually tangent circles at the vertices of the triangle, their radii are s_a, s_b , and s_c . However, they recur in various other triangle formulae: in Heron's formula for the triangle area, or as trilinear coordinates for the Mittenpunkt or barycentric coordinates for the Nagel point of the triangle. They

¹ Our answer is different than the one Woolhouse arrived at! See [31].

² Why not perimeter 1? The theory is the same either way, but if we make that choice, there will be many messy denominators to keep track of later on.

have a number of neat properties; for instance, the semiperimeter of the triangle is $s_a + s_b + s_c = 1/2(a + b + c)$, so we can again restrict to fixed perimeter by assuming that (s_a, s_b, s_c) lies on the plane $x + y + z = 1$. But now the space of all triangles is the entire orthant $s_a \geq 0, s_b \geq 0, s_c \geq 0$.

We now introduce our last set of variables: we can parametrize fixed perimeter triangle space by the unit sphere $x^2 + y^2 + z^2 = 1$ if we let

$$x^2 = s_a, \quad y^2 = s_b, \quad z^2 = s_c.$$

This is actually an eightfold cover of triangle space, but that won't make any difference to our calculations in probability. We can now solve for a, b , and c from the equations above.

Definition 1. The *symmetric measure* μ on the space of perimeter 2 triangles is given by the pushforward of the uniform probability measure on the unit sphere under the map

$$a = 1 - x^2, \quad b = 1 - y^2, \quad c = 1 - z^2.$$

The variables x, y , and z appear in various places in the theory of the triangle. We leave to the reader the (pleasant) proof of the following proposition:

Proposition 2. *Various standard quantities in triangle geometry have natural expressions in terms of the coordinates x, y , and z . For triangles with unit semiperimeter,*

- *The inradius r and triangle area A are both $|xyz|$.*
- *The three exradii r_1, r_2 , and r_3 are $|\frac{xy}{z}|, |\frac{yz}{x}|$ and $|\frac{xz}{y}|$.*
- *The variables $|x|, |y|$ and $|z|$ are the (pairwise) geometric means of the exradii.*

From this, it is easy to prove, for instance, the appealing triangle geometry theorem that $rr_1r_2r_3 = A^2$.

We have now defined a measure on triangle space, and it's clear from our construction that rotations of the sphere provide a beautiful, compact, transitive group of symmetries of triangles. This is already appealing, but one can immediately see that we have made various choices in the construction, and it is not clear how this construction would generalize to polygons with more edges. So now we start over and give another derivation of the same measure from a different point of view; this construction of polygon space is the one in our paper [3] and is originally due to Knutson and Hausmann [12].

We will start by thinking of \mathbb{R}^2 as the complex plane. Since we are interested in polygons up to translation³, we will represent a polygon by edges e_1, \dots, e_n , which are the complex numbers

³ We will deal with rotations shortly.

corresponding to the edge vectors. To fix perimeter, we want $|e_1| + \dots + |e_n| = 2$, but as before, we suspect that we will have more symmetries if we use the variables $z_1^2 = e_1, \dots, z_n^2 = e_n$ instead.

Now the polygon must close, so we are also imposing the condition $\sum e_i = 0$. If we write $z_i = u_i + v_i \mathbf{i}$, this condition becomes

$$0 = \sum z_i^2 = \sum (u_i^2 - v_i^2) + 2u_i v_i \mathbf{i}. \quad (1)$$

Rearranging, this is equivalent to $\sum u_i^2 = \sum v_i^2$ and $\sum u_i v_i = 0$. We have proved

Proposition 3. *Suppose $z_i = u_i + v_i \mathbf{i}$. The polygon with edges z_1^2, \dots, z_n^2 is closed and has perimeter 2 $\iff \vec{u} = (u_1, \dots, u_n), \vec{v} = (v_1, \dots, v_n)$ are orthonormal vectors in \mathbb{R}^n .*

Since squaring takes the 2^n points $(\pm z_1, \dots, \pm z_n)$ to the same edge set (z_1^2, \dots, z_n^2) , the Stiefel manifold of $V_2(\mathbb{R}^n)$ of orthonormal pairs of vectors in \mathbb{R}^n is a 2^n -fold cover of the space of polygons (up to translation) in the plane.⁴

An easy computation shows that rotating (\vec{u}, \vec{v}) in the plane they span rotates the polygon (twice as fast) in the plane. This means that all pairs (\vec{u}, \vec{v}) in the same plane give the same n -gon up to rotation, and so the space of n -gons up to translation and rotation is covered by the space of 2-planes in \mathbb{R}^n , which is the Grassmann manifold $G_2(\mathbb{R}^n)$.

Definition 4. The *symmetric measure* μ on the space of n -gons with perimeter 2 is given by the pushforward of the uniform probability measure on the Grassmannian $G_2(\mathbb{R}^n)$ to polygon space under the map $P \mapsto e_i = (u_i + v_i \mathbf{i})^2$, where \vec{u}, \vec{v} are any orthonormal basis for the plane P .

We have now given two definitions of the symmetric measure on triangle space, and must show they are the same.

Proposition 5. *The measures on triangle space of Definition 1 and Definition 4 are the same.*

Proof. We can identify $G_2(\mathbb{R}^3)$ with $G_1(\mathbb{R}^3)$ by taking normals to the planes; but $G_1(\mathbb{R}^3) = \mathbb{R}P^2$, which is double covered by S^2 . The uniform measure on the Grassmannian is the pushforward of the uniform measure on this sphere. This means that both measures push forward from the standard measure on the sphere; we just need to check that a given point on the sphere maps to the same triangle under each construction.

If we take a point $\vec{p} = (x, y, z)$ on the sphere, the corresponding triangle is obtained using Definition 4 by completing \vec{p} to a positive determinant orthonormal basis for \mathbb{R}^3 by adding the vectors \vec{u} and \vec{v} . Any such choices of \vec{u} and \vec{v} will produce the same triangle shape since different choices will be related by a rotation of the triangle.

Given \vec{u} and \vec{v} , the three vectors \vec{p}, \vec{u} , and \vec{v} are the columns of an orthogonal matrix. The norm of each row of the matrix is 1, so for example $u_1^2 + v_1^2 = 1 - x^2$. But $u_1^2 + v_1^2 = |(u_1 + v_1 \mathbf{i})^2| = |e_1| = a$ is the length of the first side of the triangle, so we have shown that $a = 1 - x^2$. Similarly, $b = 1 - y^2$ and $c = 1 - z^2$, as they should be according to Definition 1. \square

⁴ If one of the z_i is zero, then $z_i = -z_i$ and the order of the cover is a lower power of 2, so this cover is actually branched over the points where some $z_i = 0$, which correspond to the polygons with i th edge of length 0.

Definition 6. When sidelength $c = 1 - z^2 \neq 0$, we can define a *canonical triangle* associated to $\vec{p} = (x, y, z)$ by choosing our basis for the plane perpendicular to \vec{p} to be

$$\begin{pmatrix} \vec{u} \\ \vec{v} \end{pmatrix} = \frac{1}{\sqrt{1-z^2}} \begin{pmatrix} xz & yz & -x^2 - y^2 \\ -y & x & 0 \end{pmatrix} \quad (2)$$

It is easy to check that $\vec{p}, \vec{u}, \vec{v}$ is a positively oriented orthonormal basis for \mathbb{R}^3 . Continuing to unwind our definitions, if we place the center of edge c at the origin, then edge $c = e_3$ points in the positive x -direction, and the vertices of the triangle are

$$-\frac{1}{2}(1-z^2, 0), \quad \frac{1}{2}(1-z^2, 0), \quad \left(\frac{-x^4 + 2x^2 - 2y^2 + y^4}{2(1-z^2)}, -\frac{2xyz}{1-z^2} \right). \quad (3)$$

This observation is helpful when performing computer experiments.

2. A TRANSITIVE GROUP OF ISOMETRIES ON POLYGON SPACE

We noted above that the Grassmannian $G_2(\mathbb{R}^n)$ has a uniform measure; this is simply the unique probability measure on $G_2(\mathbb{R}^n)$ which is invariant under the left action of $O(n)$ on $G_2(\mathbb{R}^n)$. But what does this action look like on triangle space? We can start by describing the action of $SO(3)$ (rotations) on the sphere of triangles in our coordinates above: rotating around a coordinate axis fixes a sidelength, and must therefore move the opposite vertex around an ellipse as the perimeter of the triangle is fixed:

Proposition 7. *The circle $(\sqrt{1-z^2} \cos \theta, \sqrt{1-z^2} \sin \theta, z)$ formed by rotating a point on S^2 around the z -axis maps to a family of canonical triangles where vertices A and B are fixed and vertex C follows the ellipse*

$$C(\theta) = \left(\frac{1+z^2}{2} \cos 2\theta, -z \sin 2\theta \right). \quad (4)$$

Note that the ellipse is parametrized *clockwise* if $z > 0$ and counterclockwise only if $z < 0$, and that the map double-covers the ellipse. Moreover, this is the equal-area-in-equal-time parametrization of the ellipse, attesting to the naturality of this construction. The proof of the proposition is a pleasant exercise in plugging the parametrization of the circle into (3) and simplifying.

There is a rather interesting open question here: what characterizes the family of triangles obtained by an arbitrary rotation of the sphere? Infinitesimal rotations of the sphere are linear combinations of coordinate axis rotations, so we know that the family is given by integrating infinitesimal linear combinations of the above elliptical vertex motions. The resulting pictures are certainly pretty (we show an example in Figure 1), but we do not yet have a fully triangle-theoretic description of this family.

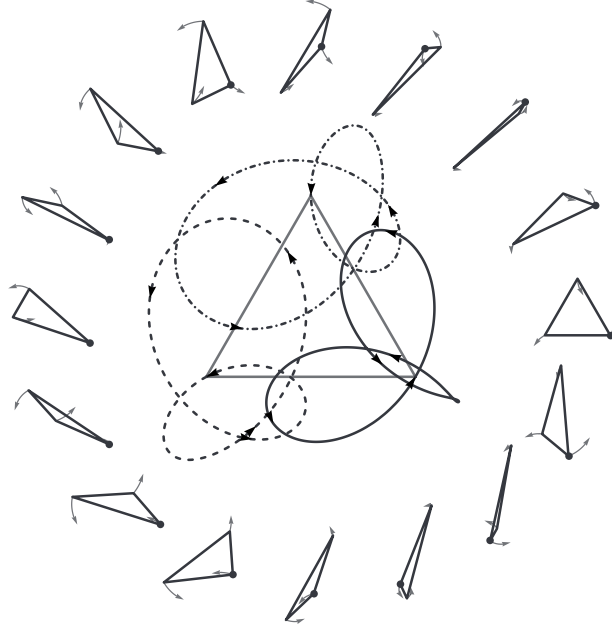


FIG. 1: These are two different visualizations of the triangle motion induced by rotating the point $\frac{1}{\sqrt{3}}(1, 1, 1)$ corresponding to the equilateral triangle around the axis $(-1, 1, -\sqrt{2})$. The circle of triangles shows 16 equally-spaced points along the resulting great circle together with the path each vertex will traverse in the next time step. The figure in the middle shows the starting equilateral triangle along with the three curves traced out by the vertices. The solid curve is the path of the vertex marked with a dot in the outside triangles.

3. THE PILLOW PROBLEM

We are now in a position to answer Lewis Carroll's pillow problem, which boils down to identifying the obtuse triangles as a subset of the unit sphere and computing its area. Recall that a point (x, y, z) on the unit sphere maps to the triangle with sidelengths

$$a = 1 - x^2, \quad b = 1 - y^2, \quad c = 1 - z^2.$$

The sphere can be split up into two (disconnected) regions: the acute triangles and the obtuse triangles; of course, the right triangles are the boundary between regions. But right triangles are easy to identify from the sidelengths: they are exactly the triangles such that $a^2 + b^2 = c^2$ or $b^2 + c^2 = a^2$ or $c^2 + a^2 = b^2$. Substituting in the above expressions for a, b, c yields three quartics:

$$(1 - x^2)^2 + (1 - y^2)^2 = (1 - z^2)^2 \tag{5}$$

and the two other cyclic permutations of the variables. The intersections of these quartics with the sphere give the collection of curves shown in Figure 2.

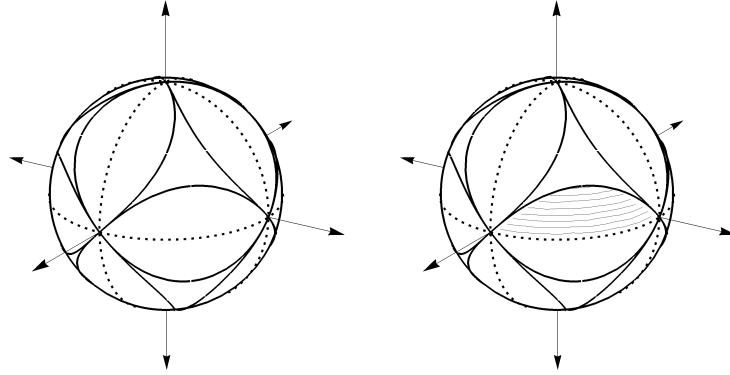


FIG. 2: The right triangles are the heavy black curves on the sphere (the dotted lines indicate the intersections of the sphere with the coordinate planes). The hatched region in the right hand figure shows $1/24$ of the region of obtuse triangles. We compute the area of this region below.

The sphere equation $x^2 + y^2 + z^2 = 1$ implies $z^2 = 1 - x^2 - y^2$, so (5) can be re-written as

$$x^2 + x^2 y^2 + y^2 = 1.$$

Equivalently, $x^2 = \frac{1-y^2}{1+y^2}$, which can be plugged into $z^2 = 1 - x^2 - y^2$ to get the following parametrization for solutions of (5):

$$\left(\pm \sqrt{\frac{1-y^2}{1+y^2}}, \pm y, \pm y \sqrt{\frac{1-y^2}{1+y^2}} \right). \quad (6)$$

Computing in cylindrical coordinates, the area of the set of obtuse triangles is simply $24 \iint_R dz d\theta$, where R is the hatched region shown at right in Figure 2.

In turn, using Stokes' Theorem,

$$24 \iint_R dz d\theta = 24 \int_{\partial R} z d\theta = 24 \left(\int_{z=0} z d\theta + \int_C z d\theta \right),$$

where C is the upper boundary of the region parametrized by (6) with all signs positive. Of course, the first integral vanishes, so we are reduced to computing the second integral. Using (6) in conjunction with $\theta = \arctan(y/x)$ to simplify yields

$$24 \int_0^1 \left(\frac{2y}{1+y^4} - \frac{y}{1+y^2} \right) dy.$$

Both terms are easy to integrate using u -substitutions: the first is recognizably the derivative of $\arctan(y^2)$, while the second is the derivative of $-\frac{1}{2} \ln(1+y^2)$, so the area of the obtuse triangles is

$$24 \left[\arctan(y^2) - \frac{1}{2} \ln(1+y^2) \right]_0^1 = 6\pi - 12 \ln 2.$$

Dividing by the area 4π of the sphere reveals the fraction of obtuse triangles to be exactly

$$\frac{3}{2} - \frac{3 \ln 2}{\pi} \approx 0.838093.$$

4. DIRICHLET DISTRIBUTIONS AND EXPECTED AREAS OF TRIANGLES

Now that we have solved Carroll's problem, it's interesting to see what other expectations we can compute! Writing things in terms of the variables $s_a = \frac{-a+b+c}{2}$, s_b and s_c , for instance, leads us to a really nice computation of the expectation of inradius (or area) and circumcurvature. Given a point (x, y, z) on the unit sphere, the corresponding triangle has $s_a = x^2$, $s_b = y^2$, and $s_c = z^2$. Then $\varphi : (x, y, z) \mapsto (x^2, y^2, z^2)$ gives a map from the unit sphere S^2 to the simplex $\{s_a + s_b + s_c = 1\}$.

Proposition 8. *The pushforward by φ of the uniform measure on S^2 is the Dirichlet distribution on the simplex with concentration parameters $(1/2, 1/2, 1/2)$. This measure is $\frac{1}{\text{Area}(\Delta ABC)} dx dy$.*

Proof. We will use x and y as coordinates on the simplex. In these coordinates, the simplex is parametrized by the triangle $x + y \leq 1$, $x \geq 0$, $y \geq 0$. The density of the uniform probability measure on S^2 (with respect to the standard area form), is just the constant function $\frac{1}{4\pi}$. But then, since φ is an 8-to-1 map, the change-of-variables formula tells us that the density of the pushforward measure on the simplex Δ is

$$8 \frac{1}{4\pi} \frac{1}{|J\varphi|}, \quad (7)$$

where $|J\varphi|$ is the Jacobian determinant of φ .

Now, we can compute the Jacobian determinant by taking the square root of the determinant of the Gramian of the 3×2 matrix $\Phi = (\nabla\varphi_1, \nabla\varphi_2)$, where $\nabla\varphi_i$ is the intrinsic gradient in S^2 of the coordinate function φ_i . Since

$$\nabla\varphi_1 = \begin{pmatrix} 2x \\ 0 \\ 0 \end{pmatrix} - 2x \begin{pmatrix} x \\ y \\ z \end{pmatrix} \quad \text{and} \quad \nabla\varphi_2 = \begin{pmatrix} 0 \\ 2y \\ 0 \end{pmatrix} - 2y \begin{pmatrix} x \\ y \\ z \end{pmatrix},$$

it is straightforward to compute

$$|J\varphi| = \sqrt{\det \Phi^T \Phi} = 4|xyz| = 4\sqrt{s_a s_b s_c}.$$

Combining this with (7), the density of the measure on $s_a + s_b + s_c = 1$ is

$$\psi(s_a, s_b) = \frac{1}{2\pi} s_a^{-1/2} s_b^{-1/2} s_c^{-1/2}, \quad (8)$$

which is the density of the Dirichlet distribution, as claimed. Since we've fixed the semiperimeter $s = 1$ for our triangles, we have $\sqrt{s_a s_b s_c} = \sqrt{s s_a s_b s_c}$, which Heron's formula says is the area of the triangle. \square

Corollary 9. *The expected value of the area of a perimeter-2 triangle with respect to the symmetric measure is*

$$E(\text{Area}) = \frac{1}{4\pi}.$$

The expression $\sqrt{ss_as_bs_c}$ also appears in the formula for the circumradius of a triangle:

$$\frac{abc}{4\sqrt{ss_as_bs_c}}$$

While the expected value of the circumradius diverges, the expected value of its reciprocal – that is, the expected curvature of the circumcircle – does not:

Corollary 10. *With respect to the symmetric measure on triangles with perimeter 2, the expected value of the curvature of the circumcircle is*

$$E(\text{circumcurvature}) = \frac{\pi}{2}.$$

Proof. The expectation of circumcurvature is

$$\iint_{\Delta} \frac{4\sqrt{ss_as_bs_c}}{abc} \psi(s_a, s_b) \, d\text{Area}.$$

Using the definition of ψ from (8) along with $s = 1$, $s_c = s - s_a - s_b = 1 - s_a - s_b$, and $a = 1 - s_a$, $b = 1 - s_b$, and $c = 1 - s_c$, this simplifies as

$$\begin{aligned} \iint_{\Delta} \frac{4}{2\pi(1-s_a)(1-s_b)(s_a+s_b)} \, d\text{Area} \\ = \int_0^1 \int_0^{1-s_a} \frac{2}{\pi(1-s_a)(1-s_b)(s_a+s_b)} \, ds_b \, ds_a = \frac{\pi}{2}. \end{aligned}$$

□

It is great fun to compute the expectation of other natural quantities in triangle geometry, and we invite the reader to continue along these lines.

5. COORDINATES FOR n -GONS

We are now going to extend our picture to n -gons. We will start by generalizing our previous coordinates x , y , and z for triangles. Remember that the vector $\vec{p} = (x, y, z)$ was the unit normal vector to the plane in $G_2(\mathbb{R}^3)$ defining the triangle, or the cross product of two orthonormal vectors \vec{u} , \vec{v} giving a basis for the plane. Each coordinate of \vec{p} is the determinant of a 2×2 matrix of coordinates $\begin{pmatrix} u_i & v_i \\ u_j & v_j \end{pmatrix}$ from \vec{u} and \vec{v} . The length of \vec{p} depends on \vec{u} and \vec{v} , but \vec{p} will always lie on the line normal to the plane. For that reason, it is useful to think of \vec{p} as defined up to scalar multiplication⁵. In \mathbb{R}^3 , there are precisely $\binom{3}{2} = 3$ such determinants, but in \mathbb{R}^n , there are $\binom{n}{2}$ such

⁵ the scalar is determined by $\det A^T A = 1/2 \sum \Delta_{ij}^2$, where A is the $n \times 2$ matrix with columns \vec{u} , \vec{v} .

determinants. This leads you to construct

Definition 11. The *Plücker coordinates* on $G_2(\mathbb{R}^n)$ associated to the plane P spanned by \vec{u} and \vec{v} are the skew-symmetric matrix $\Delta(P)$ given by taking all the 2×2 minor determinants of the $n \times 2$ matrix $A = (\vec{u} \vec{v})$ and identifying matrices which are scalar multiples of each other. This has several (immediately) equivalent forms:

$$\Delta(P)_{ij} = \det \begin{pmatrix} u_i & v_i \\ u_j & v_j \end{pmatrix} = (u_i, v_i) \times (u_j, v_j) = \left(A \begin{pmatrix} 0 & 1 \\ -1 & 0 \end{pmatrix} A^T \right)_{ij}. \quad (9)$$

When it's clear which plane the coordinates refer to, we'll just write Δ and Δ_{ij} .

We note in passing that the Plücker matrix is skew-symmetric. The matrix $\begin{pmatrix} 0 & 1 \\ -1 & 0 \end{pmatrix}$ should be familiar here; it represents multiplication by $-i$ in the standard matrix representation of complex numbers. In context, it defines a complex structure J on P given by $J(\vec{u}) = -\vec{v}$ and $J(\vec{v}) = \vec{u}$, which is to say J rotates P by 90° from \vec{v} to \vec{u} . The Plücker matrix $\Delta(P)$ now has a clear geometric interpretation: as a linear map $\Delta(P)$ orthogonally projects each $\vec{x} \in \mathbb{R}^n$ to P and then twists by the complex structure J .

The definition tells us how to find the Plücker matrix from the basis \vec{u}, \vec{v} , and this geometric interpretation (or the last expression in (9)) tell us how to go back. Since $\Delta(P)\vec{u} = -\vec{v}$ and $\Delta(P)\vec{v} = \vec{u}$, the pairs $(\vec{u}, -\vec{v})$ and (\vec{v}, \vec{u}) are singular vector pairs associated to the singular value 1 for $\Delta(P)$. This means that the two left singular vectors or the two right singular vectors corresponding to the singular value 1 also give an orthonormal basis for the plane. Hence, we can recover an orthonormal basis for the plane from the Plücker coordinates by taking the SVD of the matrix $\Delta(P)$. Since the singular values are known in advance (two singular values are 1, the rest are 0), this is constructive and exact.

Last, we remark that while every Plücker matrix is skew-symmetric, the Plücker matrices are only a small subset of the skew-symmetric matrices: the Plücker coordinates obey an interesting system of Plücker relations which encode the fact that these subdeterminants are not all independent. The super-diagonal entries of the Plücker matrix define homogeneous coordinates of an embedding of $G_2(\mathbb{R}^n)$ into the projective space $\mathbb{RP}^{\binom{n}{2}-1}$. Since $\dim G_2(\mathbb{R}^n) = 2(n-2)$ is less than $\dim \mathbb{RP}^{\binom{n}{2}-1} = \binom{n}{2} - 1$ for $n \geq 4$, we expect that these coordinates satisfy additional constraints. In fact, the constraints are simple: for each choice of four distinct rows $i < j < k < \ell$ from the matrix $A = (\vec{u} \vec{v})$, there are six Plücker coordinates $\Delta_{ij}, \Delta_{ik}, \Delta_{il}, \Delta_{jk}, \Delta_{j\ell}, \Delta_{k\ell}$ coming from the six possible 2×2 minors involving the four rows. These six coordinates must satisfy the *Plücker relation*

$$\Delta_{ij}\Delta_{k\ell} - \Delta_{ik}\Delta_{j\ell} + \Delta_{il}\Delta_{jk} = 0. \quad (10)$$

The relations coming from all possible choices of four rows define a system of homogeneous quadratic equations which exactly cut out the image of the Grassmannian inside projective space. See [17] for a beautifully clear discussion of these matters.

6. FROM CONTINUOUS SYMMETRIES TO DISCRETE SYMMETRIES

We have now shown that $O(3)$ is a transitive group which acts on triangles and preserves the measure; generalizing, Definition 4 tells us that $O(n)$ does the same for n -gons. This is the start of a fascinating journey, as the structure of the orthogonal group is one of the most beautiful chapters in algebra. Every math student knows that there are only 5 Platonic solids; more advanced ones know that the relatively scarcity of these extraordinary shapes comes from the fact there are only a few finite subgroups of $O(3)$. However, in higher dimensions there are many more finite subgroups, and each of these yields an beautiful symmetry of polygon space.

We focus on the hyperoctahedral group B_n , which is the subgroup of matrices in $O(n)$ of signed permutations of the coordinates x_1, \dots, x_n ; it is the group of symmetries of the hypercube and of its dual, the cross-polytope or hyperoctahedron. These are the only matrices in $O(n)$ with integer coordinates: each such matrix can be written as the product of a diagonal matrix with entries ± 1 and a permutation matrix. It will be most convenient to describe an arbitrary element $\beta \in B_n$ as a permutation of $(-n, \dots, -1, 1, \dots, n)$ obeying the condition $\beta(-i) = -\beta(i)$. Our goal now is to describe the action of the hyperoctahedral group on polygon space.

The action of B_n on a plane $P \in G_2(\mathbb{R}^n)$ permutes the rows of any basis (\vec{u}, \vec{v}) for P and changes some of their signs. This action is never effective: reversing the sign of all rows yields the same plane, though this element is the only nontrivial stabilizer of a generic plane. The action descends to an action on polygon space. A generic polygon is stabilized by the $(\mathbb{Z}/2\mathbb{Z})^n$ subgroup of B_n of signed permutations β where $\beta(i) = \pm i$ for all i . The quotient group of unsigned permutations $S_n = B_n/(\mathbb{Z}/2\mathbb{Z})^n$ simply permutes the edges of the polygon. These facts prove that

Proposition 12. *The symmetric measure on polygon space is invariant under permutations of the edges.*

We now want to explore some consequences of this invariance. It follows directly from the definitions that β acts in a nice way on Plücker matrices:

Proposition 13. *For any $\beta \in B_n$,*

$$\Delta(\beta P)_{ij} = \text{sgn } \beta(i) \text{sgn } \beta(j) \Delta(P)_{|\beta(i)||\beta(j)|}.$$

We now start describing subsets of polygon space:

Definition 14. We say that a 2-plane P with orthonormal basis $A = (\vec{u} \vec{v})$ is a *semicircular lift* of a polygon if the directions of the vectors (u_i, v_i) all lie on the semicircle (oriented counterclockwise) between (u_1, v_1) and $-(u_1, v_1)$.

We might worry that this idea is not well-defined; after all, there are many orthonormal bases for the plane! But it is easy to see that changing bases just rigidly rotates the collection of vectors (u_i, v_i) , which preserves the property above.

7. POLYGONS AND THE POSITIVE GRASSMANNIAN

A subset of the Grassmannian which has attracted a lot of interest recently in string theory [1] is the *positive Grassmannian* of planes with a basis for which $\Delta_{ij}(P) > 0 \iff i < j$. We note that any basis for a plane in the positive Grassmannian has all signs in the upper triangle agreeing, but that reversing the orientation of the plane (for instance) reverses the signs of all Plücker coordinates. Therefore, we might also see matrices with all negative signs in the upper triangle. (In this case we took the wrong basis.)

This subspace has a natural meaning in our terms:

Proposition 15. *The positive Grassmannian $G_2(\mathbb{R}^n)^+$ consists of planes with a basis which is a semicircular lift of a strictly convex polygon.*

The proof is a pleasant exercise in chasing down definitions, so we leave it to the reader with this hint: strict convexity of the polygon is equivalent to the statement that the edge directions are distinct and in counterclockwise order on the circle and a semicircular lift preserves this property. We note that a very similar interpretation of the positive Grassmannian for 2-planes which shows that the positive Grassmannian has the same topology as the convex polygons appears in Section 5.3 of [1].

Since the property of being a convex polygon is invariant under cyclic permutations of the edges, we expect a cyclic subgroup of the hyperoctahedral group to preserve $G_2(\mathbb{R}^n)^+$. In fact, the full stabilizer is somewhat bigger, and the cyclic part is not quite what we'd expect:

Proposition 16. *The stabilizer of $G_2(\mathbb{R}^n)^+$ inside the hyperoctahedral group is the subgroup of order $4n$ generated by*

$$\begin{aligned}\beta &= (1, 2, \dots, (n-1), -n)(-1, -2, \dots, -(n-1), n), \\ \eta &= (-1, 1)(-2, 2) \dots (-n, n), \\ \gamma &= (1, n)(2, n-1) \dots (-1, -n)(-2, -(n-1)) \dots\end{aligned}$$

Note that the subgroup generated by β is cyclic of order n , but not the canonical cyclic subgroup of order n generated by $(1, 2, \dots, n)(-1, -2, \dots, -n)$.

Proof. We start by showing that all these group elements map $G_2(\mathbb{R}^n)^+$ to itself, using Proposition 13. If P is in $G_2(\mathbb{R}^n)^+$, then

$$\Delta(\beta P)_{ij} = \Delta(P)_{(i+1)(j+1)} > 0 \iff i < j \text{ and } i, j \neq n$$

while

$$\Delta(\beta P)_{in} = -\Delta(P)_{(i+1)1} > 0, \text{ since } \Delta(P)_{1(i+1)} > 0,$$

so βP is still in $G_2(\mathbb{R}^n)^+$. The element η doesn't change any Plücker coordinates:

$$\Delta(\eta P)_{ij} = \text{sgn } \eta(i) \text{sgn } \eta(j) \Delta(P)_{ij} = (-1)(-1) \Delta(P)_{ij} = \Delta(P)_{ij}.$$

Thus ηP is actually the same plane! (And so it's definitely still in $G_2(\mathbb{R}^n)^+$.) Since γ reverses the order of the coordinates, $i > j \iff \gamma(i) < \gamma(j)$. Further, $\Delta(P)_{ij} > 0 \iff i < j$. Thus,

$$\Delta(\gamma P)_{ij} = \Delta(P)_{\gamma(i)\gamma(j)} > 0 \iff i > j.$$

Thus $\Delta(\gamma P)$ has positive entries *below* the main diagonal and negative entries *above* in some basis (\vec{u}, \vec{v}) for P . But the basis (\vec{v}, \vec{u}) for P has opposite signs for all Plücker coordinates, and so has $\Delta_{ij} > 0 \iff i < j$, as desired. Thus $\gamma P \in G_2(\mathbb{R}^n)^+$.

It's fun to see this geometrically as well. Figure 3 shows a convex 4-gon and its semicircular lift. If we cyclically permute the edges so that we start with edge 4, we must take the other square root of edge direction 4 to keep the lifts of the other edges in the semicircle extending counterclockwise from (u_4, v_4) . This generalizes to n -gons.

We have now proved that our subgroup stabilizes $G_2(\mathbb{R}^n)^+$. But we haven't proved that it is the *largest* subgroup of the hyperoctahedral group with this property. So take any hyperoctahedral group element δ which stabilizes $G_2(\mathbb{R}^n)^+$. We know δ must send 1 to some $\pm k$. The product π of δ and $\beta^{n-|k|}$ (and, if needed, η) fixes 1. It suffices to show that π is the identity; if so, δ was in the subgroup.

We start by proving that $\pi(j) > 0$ for all positive j . Suppose not. Then if P is in the positive Grassmannian,

$$\Delta(\pi P)_{1j} = \text{sgn } \pi(1) \text{sgn } \pi(j) \Delta P_{1|\pi(j)|} = -\Delta P_{1|\pi(j)|} < 0,$$

since $\Delta P_{1|\pi(j)|} > 0$ by our assumption that P was in the positive Grassmannian. Thus πP is not in the positive Grassmannian, a contradiction. This means that π consists of matching permutations of $1, \dots, n$ and $-1, \dots, -n$.

We are now going to prove by induction that $\pi(k) = k$ for all k . We have just established the base case ($k = 1$). So suppose $\pi(k) = k$ for $k < K$, and consider $\pi(K)$. If $\pi(K) \neq K$, then $\pi(K) > K$ (since $1, \dots, K-1$ are taken), and there is some $j > K$ so that $\pi(j) = K$. But then

$$\Delta(\pi P)_{Kj} = \Delta P_{\pi(K)K} < 0,$$

even though $\pi(K) > K$, and πP is not in the positive Grassmannian, a contradiction. Thus $\pi(K) = K$, and we have proved that π is the identity permutation. \square

We can now subdivide the space of polygons in a natural way by dividing the Grassmannian into *sign chambers* by grouping together all planes for which the matrix $S_{ij} = \text{sgn } \Delta_{ij}$ of signs of the Plücker coordinates is the same. We will call these *Plücker sign matrices*. By convention, the sign chambers will be the open subsets of the Grassmannian whose Plücker sign matrices have zeros only on the diagonal. The positive Grassmannian, for instance, is the sign chamber corresponding to the Plücker sign matrix S^0 defined by $S_{ij}^0 = 1 \iff i < j$, $S_{ii}^0 = 0$.

Like the original Plücker matrices, Plücker sign matrices are skew-symmetric and defined up to scalar multiplication (by ± 1). But not every skew-symmetric matrix of ± 1 's is a Plücker sign matrix – the Plücker relations (10) rule some out. This means that it is interesting to count the sign

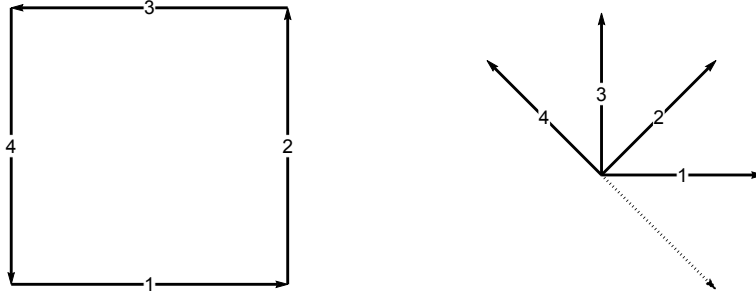


FIG. 3: A convex polygon in the positive Grassmannian $G_2(\mathbb{R}^n)^+$ (left) and its semicircular lift (right). If we cyclically permute the edges to put edge 4 first, we must take the opposite lift of edge 4 to make the lift semicircular.

chambers and determine whether there are different types of sign chambers or whether they are all identical.

Proposition 17. *The action of the hyperoctahedral group on $G_2(\mathbb{R}^n)$ descends to a transitive hyperoctahedral group action on the sign chambers and the Plücker sign matrices.*

Proof. It follows immediately from Proposition 13 that the action of the hyperoctahedral group on the Grassmannian and the Plücker matrices induces a corresponding action on the set of Plücker sign matrices. So suppose we have an arbitrary Plücker sign matrix S corresponding to some plane $P \in G_2(\mathbb{R}^n)$ with a basis (\vec{u}, \vec{v}) , as usual. It suffices to show that there is a hyperoctahedral group element which maps S to S^0 .

Some collection of sign changes puts all the $(u_i, v_i) = u_i + v_i \mathbf{i} = r_i e^{i\theta_i}$ in the semicircle extending counterclockwise from $(u_1, v_1) = (r_1, 0)$. There is then a unique permutation of $2, \dots, n$ which fixes (u_1, v_1) and puts the remaining (u_i, v_i) in counterclockwise order by direction. Together, the permutation and sign changes are some element β of the hyperoctahedral group. But now the basis \vec{u}, \vec{v} is a semicircular lift of a convex polygon and by Proposition 15, the resulting plane is in the positive Grassmannian. Hence, it has Plücker sign matrix S^0 . \square

We can now count and describe the sign chambers easily:

Proposition 18. *There are $2^{n-2} \times (n-1)!$ sign chambers and corresponding Plücker sign matrices. The sign chambers are all isometric and in particular have the same volume.*

Proof. By the orbit-stabilizer theorem, the size of the orbit of S^0 is equal to the size of the hyperoctahedral group B_n (namely $2^n \times n!$) divided by the size of the stabilizer ($4n$, by Proposition 16). But by Proposition 17, the orbit of S^0 is the entire set of Plücker sign matrices. \square

We now want to understand the geometric meaning of the sign chambers.

Theorem 19. *The convex n -gons consist of 2^{n-1} copies of the positive Grassmannian. They comprise $2/(n-1)!$ of the space of n -gons.*

Proof. We already know from Proposition 15 that the positive Grassmannian consists of the semi-circular lifts of the convex polygons. Therefore, any other plane corresponding to a convex polygon must be a different lift to $G_2(\mathbb{R}^n)$. There are 2^{n-1} such lifts, remembering that changing all the signs has no effect. \square

8. SYLVESTER'S 4-POINT PROBLEM AND QUADRILATERALS

In the same ongoing discussion in which Woolhouse posed his versions of the obtuse triangle problem, J. J. Sylvester in 1864 asked for the probability that four points “taken at random in a plane” formed the vertices of a reentrant (embedded, but not convex) quadrilateral [25].⁶ Various solutions were proposed by Cayley [24, footnote 64(b)], De Morgan [7, pp. 147–148], and others, with answers including (at least) $1/4$, $35/12\pi^2$, $3/8$, $1/3$, and $1/2$ [14]. As with Carroll’s problem, it soon became clear that the probability measure for the four points was an issue, with Sylvester concluding that the triangle problem and the quadrilateral problem, as posed, “do not admit of a determinate solution” [26]. A robust literature has grown up around the related problem of finding the probability when the points are selected from the interior of a convex body (see in particular Blaschke’s remarkable result [2] and Pfeifer’s survey [22]).

From our perspective, the most compelling of the original solutions to Sylvester’s problem was given by the science educator, astronomer, future priest, and past Senior Wrangler James Maurice Wilson [28], who argued that $1/3$ of quadrilaterals are reentrant by focusing on the edges of the quadrilateral rather than the vertices. Indeed, an extrapolation of his argument suggests that $1/3$ of quadrilaterals should be convex, $1/3$ reentrant, and $1/3$ self-intersecting, the same answer we will arrive at in Theorem 26.

We now answer Sylvester’s question in our terms. We divide quadrilaterals into three classes: convex, reflex (or reentrant), and self-intersecting. We have shown (Theorem 19) that $1/3$ of the quadrilaterals are convex; the remaining 4-gons are either reflex or self-intersecting. The boundaries between the classes consist of polygons where two edges point in the same (or opposite) directions; that is, when rows of the $n \times 2$ matrix A are colinear or perpendicular. The Plücker matrix consists of cross products of these rows, and detects colinearity; we now add the matrix of dot products of rows to detect perpendicularity:

Definition 20. The *projection matrix* associated to a plane P with (orthonormal) basis given by the $n \times 2$ matrix A is given by AA^T . This matrix orthogonally projects vectors to the plane P . The entries of AA^T are the dot products of the rows of A ; $(AA^T)_{ij} = (u_i, v_i) \cdot (u_j, v_j)$.

⁶ Note that there is a typo in the original statement of Sylvester’s problem: he used the word “convex” where he meant to say “reentrant”.

It is a neat fact that the projection matrix is closely related to the Plücker matrix!

Proposition 21. *If A is an orthogonal $n \times 2$ matrix, the projection matrix*

$$AA^T = -(\Delta(P))^2.$$

Proof. We noted in Definition 11 that $\Delta(P) = A \begin{pmatrix} 0 & 1 \\ -1 & 0 \end{pmatrix} A^T$. Expanding $\Delta(P)^2$:

$$\Delta(P)^2 = A \begin{pmatrix} 0 & 1 \\ -1 & 0 \end{pmatrix} A^T A \begin{pmatrix} 0 & 1 \\ -1 & 0 \end{pmatrix} A^T = A \begin{pmatrix} 0 & 1 \\ -1 & 0 \end{pmatrix}^2 A^T = A(-I)A^T = -AA^T,$$

using the fact that Gramian $A^T A = I$ since the columns of A are orthonormal.

Geometrically this is similarly clear: $-\Delta(P)^2$ has the effect of projecting a vector to P , rotating it by 180° , and then reversing its direction. Since the last two actions cancel each other, this is just projection to P . \square

We can use the Plücker sign matrix $\text{sgn } \Delta(P)$ and the *projection sign matrix* $\text{sgn } AA^T$ to divide the Grassmannian into natural cells. We will call these “sign cells” for now. As before, the hyperoctahedral group B_n acts on the matrices AA^T and $\text{sgn } AA^T$ just as it did on $\Delta(P)$ and $\text{sgn } \Delta(P)$, and we will use this group action to do our computations.

Definition 22. The Grassmannian is divided into a collection of cells, called *sign cells* where each P belongs to the subspace of all planes with the same matrices $\text{sgn } \Delta(P)$ of signs of Plücker coordinates and $\text{sgn } AA^T$ of signs of entries in the projection matrix for P .

We will now specialize to $G_2(\mathbb{R}^4)$ and prove some useful facts about the sign cells:

Proposition 23. • *The positive Grassmannian $G_2(\mathbb{R}^4)^+$ is divided into 4 sign cells.*

- *The stabilizer of the positive Grassmannian from Proposition 16 acts transitively on these 4 cells; hence the hyperoctahedral group acts transitively on the sign cells of $G_2(\mathbb{R}^4)$.*
- *The stabilizer of the “base” sign cell*

$$\text{sgn } \Delta = \begin{pmatrix} 0 & 1 & 1 & 1 \\ -1 & 0 & 1 & 1 \\ -1 & -1 & 0 & 1 \\ -1 & -1 & -1 & 0 \end{pmatrix} \quad \text{sgn } AA^T = \begin{pmatrix} 1 & 1 & 1 & -1 \\ 1 & 1 & 1 & 1 \\ 1 & 1 & 1 & 1 \\ -1 & 1 & 1 & 1 \end{pmatrix}$$

consists only of the 4 element group generated by

$$\begin{aligned} \eta &= (-1, 1)(-2, 2) \dots (-n, n) \\ \gamma &= (1, n)(2, n-1) \dots (-1, -n), (-2, -(n-1)), \dots \end{aligned}$$

- *There are $384 = 2^4 \times 4!$ elements of the hyperoctahedral group B_4 and hence $96 = 384/4$ different sign cells.*

- Each sign cell is equiprobable (in fact, each is isometric).

Proof. Using Proposition 21 and the fact that $\Delta_{ij} = -\Delta_{ji}$, we can write the projection matrix AA^T in terms of the Plücker coordinates. Since AA^T is symmetric and since the diagonal entries are simply the squared norms of the rows of Δ and hence positive, the projection sign matrix is completely determined by the super-diagonal triangle of AA^T , which is

$$\begin{pmatrix} * & \Delta_{13}\Delta_{23} + \Delta_{14}\Delta_{24} & \Delta_{14}\Delta_{34} - \Delta_{12}\Delta_{23} & -\Delta_{12}\Delta_{24} - \Delta_{13}\Delta_{34} \\ * & * & \Delta_{12}\Delta_{13} + \Delta_{24}\Delta_{34} & \Delta_{12}\Delta_{14} - \Delta_{23}\Delta_{34} \\ * & * & * & \Delta_{13}\Delta_{14} + \Delta_{23}\Delta_{24} \\ * & * & * & * \end{pmatrix} \quad (11)$$

(the elements replaced with * don't affect our analysis). For any element of the positive Grassmannian we know that $\Delta_{ij} > 0$ for all $i < j$, so the only entries in AA^T whose signs could disagree with the base sign cell are $(AA^T)_{13}$ and $(AA^T)_{24}$. Since there are only 4 possible sign combinations for these two entries, we see that there are no more than 4 sign cells in the positive Grassmannian.

Now, we examine the action of the stabilizer of the positive Grassmannian. As we saw in Proposition 16, the element η does not affect the Plücker coordinates, and hence must also fix the projection sign matrix. By inspection of the above expression for AA^T , the element γ reflects the entries of the projection matrix across the anti-diagonal, which also has no effect on the projection sign matrix of elements of the base sign cell.

The action of β is more complicated, though straightforward enough to write down using our expression for AA^T . Among other things, β replaces $(AA^T)_{13}$ with $(AA^T)_{24}$ and replaces $(AA^T)_{24}$ with $-(AA^T)_{13}$. When applied to the base sign cell, this turns out to be the only way in which β affects the projection sign matrix. In turn, this means that β^2 replaces $(AA^T)_{13}$ with $-(AA^T)_{13}$ and replaces $(AA^T)_{24}$ with $-(AA^T)_{24}$, and that β^3 replaces $(AA^T)_{13}$ with $(AA^T)_{24}$ and replaces $(AA^T)_{24}$ with $(AA^T)_{13}$. In particular, all 4 possible sign cells are actually realized and the cyclic subgroup generated by β acts freely and transitively on them.

Since the stabilizer of the base sign cell must be contained in the stabilizer of the entire positive Grassmannian, we have shown that the stabilizer is exactly the subgroup generated by η and γ . The count of sign cells and the proof that they are isometric now go as they did in the proof of Proposition 17. \square

We can now show:

Proposition 24. *The sign cells in $G_2(\mathbb{R}^4)$ are path-connected.*

Proof. By Proposition 23, it suffices to prove this for the base sign cell since all the sign cells are isometric. So suppose we have two planes P_0 and P_1 in the base sign cell, with Plücker coordinates Δ_{ij}^0 and Δ_{ij}^1 . Keeping in mind that $\Delta_{ij} = -\Delta_{ji}$, we can restrict our attention to the six coordinates Δ_{ij} with $i < j$ for the remainder of the proof. (The complementary Plücker coordinates follow a

similar argument with various signs and inequalities reversed, which we will not write out.) These are all positive numbers which must obey the single Plücker relation

$$\Delta_{12}\Delta_{34} - \Delta_{13}\Delta_{24} + \Delta_{14}\Delta_{23} = 0. \quad (12)$$

Our strategy is to interpolate between P_0 and P_1 . While there is certainly a geodesic path joining P_0 and P_1 which would be a natural candidate for interpolation, that path does not seem to always stay within the sign cell. Thus, we will join P_0 and P_1 by interpolating between Δ_{ij}^0 and Δ_{ij}^1 by a family of Plücker coordinates $\Delta_{ij}(t)$.

We define the $\Delta_{ij}(t)$ by the logarithmic interpolation⁷

$$\Delta_{ij}(t) = (\Delta_{ij}^0)^{1-t} (\Delta_{ij}^1)^t \quad (13)$$

except for $\Delta_{24}(t)$, which must be given by

$$\Delta_{24}(t) = \frac{\Delta_{12}(t)\Delta_{34}(t) + \Delta_{14}(t)\Delta_{23}(t)}{\Delta_{13}(t)}$$

in order to ensure that the $\Delta_{ij}(t)$ obey the Plücker relation (12). Since the Δ_{ij}^0 and Δ_{ij}^1 are positive, it is straightforward to check that all the $\Delta_{ij}(t)$ are also positive.

We must now prove that the $(AA^T)_{ij}(t)$ have the correct signs. As in (11), we can write each $(AA^T)_{ij}(t)$ in terms of the $\Delta_{ij}(t)$ and, since $\Delta_{ij}(t) > 0$, only

$$\begin{aligned} (AA^T)_{13}(t) &= -\Delta_{12}(t)\Delta_{23}(t) + \Delta_{14}(t)\Delta_{34}(t) \text{ and} \\ (AA^T)_{24}(t) &= \Delta_{12}(t)\Delta_{14}(t) - \Delta_{23}(t)\Delta_{34}(t) \end{aligned}$$

could change sign. So it suffices to show $(AA^T)_{13}(t) > 0$ and $(AA^T)_{24}(t) > 0$ knowing that these inequalities are satisfied for $t = 0$ and $t = 1$. Rearranging, this is equivalent to showing that for all t ,

$$\frac{\Delta_{14}(t)}{\Delta_{23}(t)} > \frac{\Delta_{12}(t)}{\Delta_{34}(t)} \quad \text{and} \quad \frac{\Delta_{14}(t)}{\Delta_{23}(t)} > \frac{\Delta_{34}(t)}{\Delta_{12}(t)}. \quad (14)$$

The form of these inequalities explains why we chose logarithmic interpolation rather than linear interpolation: linearly interpolating the numerator and denominator of a fraction has a complicated effect on the quotient, while logarithmically interpolating the numerator and denominator logarithmically interpolates their quotient. So

$$\frac{\Delta_{ij}(t)}{\Delta_{kl}(t)} = \frac{(\Delta_{ij}^0)^{1-t} (\Delta_{ij}^1)^t}{(\Delta_{kl}^0)^{1-t} (\Delta_{kl}^1)^t} = \left(\frac{\Delta_{ij}^0}{\Delta_{kl}^0} \right)^{1-t} \left(\frac{\Delta_{ij}^1}{\Delta_{kl}^1} \right)^t.$$

⁷ So called because $\ln \Delta_{ij}(t)$ interpolates linearly between $\ln \Delta_{ij}^0$ and $\ln \Delta_{ij}^1$.

This means that both sides of the inequalities in (14) are being logarithmically interpolated between values at $t = 0$ and $t = 1$ where the inequalities are obeyed. It is a general fact about logarithmic interpolation that this implies the inequalities are satisfied for intermediate values of t . Proving this is a fun exercise: start by taking the logarithm of the inequalities in (14). \square

We can now see that the sign cells have geometric meaning at the level of quadrilaterals:

Proposition 25. *All of the quadrilaterals in any given sign cell are either convex, reflex, or self-intersecting.*

Proof. The walls between sign cells consist of planes with a Plücker coordinate or entry in the projection matrix equal to zero. Recall that the edges of the polygon e_i are equal to the squares of the complex numbers $z_i = u_i + \mathbf{i}v_i$.

When the Plücker coordinate $\Delta_{ij} = 0$, two rows (u_i, v_i) and (u_j, v_j) are colinear, and $z_i = \lambda z_j$ for some real λ . Squaring, we see that $e_i = \lambda^2 e_j$, and the edges e_i and e_j point in the same direction.

When the projection coordinate $(AA^T)_{ij} = 0$, two rows (u_i, v_i) and (u_j, v_j) are perpendicular, and $z_i = \lambda \mathbf{i} z_j$ for some real λ . Squaring, we see that now $e_i = -\lambda^2 e_j$, and the edges e_i and e_j point in opposite directions.

Since we have showed that the sign cells are path-connected in Proposition 24, the polygons in each sign cell can be deformed to one another without making any pair of edge directions parallel or antiparallel. But to transition between convex, reflex, and self-intersecting, two adjacent edges must point in the same or opposite directions. \square

We pause for a minute to appreciate the significance of this result. We now know that $G_2(\mathbb{R}^4)$ is broken up into 96 isometric pieces, each of which corresponds to a collection of quadrilaterals which are all convex, all reflex, or all self-intersecting. This means that we can determine which category *all* elements of a given sign cell are in by choosing a *particular* element of the sign cell and determining whether it is convex, reflex, or self-intersecting.

Since each of the sign cells is the image of the base sign cell under the action of 4 different elements of the hyperoctahedral group B_4 (namely, conjugates of the stabilizer of the base sign cell), we can choose this preferred element of each sign cell by choosing a preferred element of the base sign cell and then moving it around by hyperoctahedral elements.

In other words, to see what fraction of quadrilaterals are convex, reflex, and self-intersecting, we simply need to look at the 96 images of the special quadrilateral in the base sign cell and determine which fraction fall into each category. The answer is simple and pleasant:

Theorem 26. *Of the 96 sign cells in $G_2(\mathbb{R}^4)$, convex, reflex, and self-intersecting quadrilaterals each compose 32 cells. Hence, the probabilities that a randomly selected quadrilateral is convex, reflex, or self-intersecting are each equal to 1/3.*

Proof. As discussed above, it suffices to choose an element of the base sign cell and examine its images under the action of the hyperoctahedral group. While this will produce points in each of

the 96 different sign cells, many of the resulting quadrilaterals will be identical: after all, $B_4 \simeq (\mathbb{Z}/2)^4 \rtimes S_4$, but the normal subgroup $(\mathbb{Z}/2)^4$ corresponds to different lifts of the quadrilateral, which give different points on the Grassmannian that map to the same quadrilateral.

Hence, we can simply act on our chosen element of the base sign cell by the quotient S_4 , which just acts by permuting the edges. Since each of the 96 quadrilaterals we are tabulating is identical to one of these 24 permutation images, the theorem will follow if we see that 8 are convex, 8 are reflex, and 8 are self-intersecting.

Since we can generate random points in $G_2(\mathbb{R}^4)$, it is easy to pick a random element of the base sign cell. But the most beautiful such choice would be a “center of mass” or “average” plane in the sign cell.

(We pause for a moment to reflect on the fact that before we started, the project of defining the average of a collection of quadrilaterals would have been rather daunting. But since we are working on the Grassmannian, we have a powerful collection of tools adapted to exactly these problems!)

One definition of an average of a finite collection of subspaces of \mathbb{R}^n is given by the *flag mean* of the points. Given $n \times 2$ orthonormal matrices A_1, \dots, A_m giving bases for the subspaces, a basis for the 2-dimensional flag mean subspace is given by the two (left) singular vectors of the $n \times 2m$ matrix $(A_1 \dots A_m)$. The flag mean is used in signal processing, and has beautiful mathematical properties; see [19] and [8].

We computed the flag mean of 10,420 points sampled uniformly from the base sign cell to be the quadrilateral with edge vectors approximately $(0.33, -0.59)$, $(-0.29, -0.13)$, $(-0.30, 0.11)$, and $(0.26, 0.62)$. This quadrilateral and its 23 companions given by permuting the edges form a collection of ideal quadrilaterals representing each of the sign cells. We complete the proof by presenting these quadrilaterals in Figure 4. The reader can easily verify that 8 are convex, 8 are reflex, and 8 are self-intersecting. \square

Figure 4 is a bit disappointing at first: despite our promise of 24 different quadrilaterals there are really only 3 that are geometrically different, one in each class. But this is easily explained: cyclically permuting the edges of a quadrilateral or reversing their order produces an ordered quadrilateral which is congruent to the original, so the standard copy of the dihedral group D_8 inside S_4 produces congruent quadrilaterals. The fact that we see 3 geometrically distinct quadrilaterals thus boils down to the fact that $|S_4|/|D_8| = 3$.

We leave it to the reader to check that the order 3 subgroup $A_3 = \{1, (123), (132)\}$ is a transversal of the dihedral group D_8 inside S_4 , and hence that the three geometrically distinct quadrilaterals in Figure 4 can be obtained by applying these three permutations to our chosen representative of the base sign cell.

We have now given our solution to Sylvester’s problem. But notice that we have actually done much more: we have given an explicit geometry to the space of unordered (length 2) quadrilaterals by identifying the convex, reflex, and self-intersecting quadrilaterals as the three isometric Riemannian manifolds with boundary comprising the A_3 -orbit of the base sign cell, each a Riemannian submanifold of $G_2(\mathbb{R}^4)$.

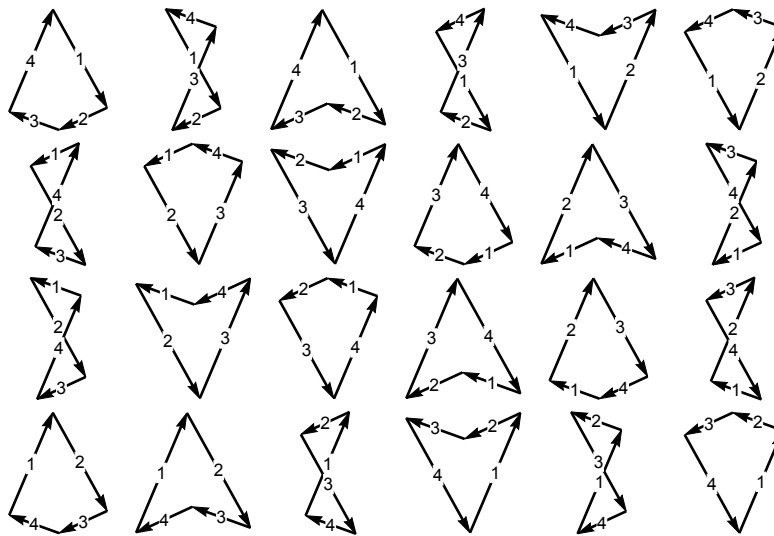


FIG. 4: The permutation group orbit of the quadrilateral corresponding to the flag mean of the base sign cell. The flag mean of every sign cell corresponds to a quadrilateral congruent to one of these 24 ordered quadrilateral. By inspection, $1/3$ are convex, $1/3$ are reflex, and $1/3$ are self-intersecting.

9. CONCLUSION AND OPEN QUESTIONS

Our journey has taken us through some beautifully concrete applications of the Grassmannian picture of planar polygons, but there is still a vast landscape to explore. Even for triangles, there are still a number of interesting questions open. First, it would be really interesting to be able to characterize the distance between triangles (and the effect of an arbitrary rotation of the sphere) directly in terms of triangle geometry. It is clear that you can write a rotation of the sphere in terms of a conserved quantity which you can express in terms of sidelengths of the triangle (we leave the exercise to the reader). But what does this formula mean? Second, it is tempting to go back and reprove many of the standard algebraic identities connecting various measurements of the triangle in terms of these variables (see, for instance, [18] for a trove of such identities concerning the inradii and exradii).

For quadrilaterals, a number of interesting open questions remain. The space $G_2(\mathbb{R}^4)$ has an interesting involution: take each plane to the perpendicular one. This gives rise to an involution on quadrilaterals which seems fascinating to explore. Many of our statements about the structure of $G_2(\mathbb{R}^4)$ should generalize to $G_2(\mathbb{R}^n)$. For instance, it seems clear that the sign cells of $G_2(\mathbb{R}^n)$ should be path-connected. Can you prove it?

This metric on plane polygons can be extended to a corresponding metric on plane curves, which is used in shape recognition and classification. The paper of Younes et al. [32] is a great place to start reading about this topic. The flag mean, and other tools from signal processing, also seem to have fruitful applications in polygon space. For starters: can you define the flag mean of

an arbitrary subset of the Grassmannian by integration rigorously? If so, is the flag mean of the base sign cell the kite we show above?

Several authors, notably Hausmann and Knutson [12], Kapovich and Millson [15], and Howard et al. [13], have extended this structure to space polygons, and specialized it to polygons with fixed edgelengths. In our own papers [3–6, 20, 21] we have developed the sampling and integration theory for these polygon spaces and extended the theory to space curves. Space polygons of fixed edgelength form a symplectic manifold, which begins another long and fascinating story. Interestingly, that manifold seems to genuinely have more structure than the space of *planar* polygons with fixed edgelengths, which remains somewhat mysterious. We hope to address fixed edgelength planar polygons in more detail in a future publication.

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