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## Serret-Frenet equations

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§1. A curve in $\mathbb{R}^{d}$ is a map from a subset of $\mathbb{R}$ into $\mathbb{R}^{d}$. A curve is called immersed is it is differentiable with everywhere nonzero derivative. The local differential geometry of curves $\alpha$ is about the properties of the set

$$
\operatorname{image}(\alpha):=\{\alpha(t) \mid t \in \operatorname{domain}(\alpha)\}
$$

which depend on the values of $\alpha$ on arbitrarily small open intervals of the domain. The term geometry is used because these properties depend on the image set, as opposed to the function $\alpha(t)$. For example,

$$
\alpha(t):=(\cos t, \sin t, 0)
$$

is a $C^{\infty}$ immersed curve in $\mathbb{R}^{3}$, and image $(\alpha)$ is the circle of radius 1 centered at the origin. The study of circles is properly described as geometry. The curve

$$
\beta(t):=(\cos 2 t, \sin 2 t, 0)
$$

is a different function, but $\alpha$ and $\beta$ have the the same image and the same local geometry.
$\S 2$. The motion of a car along a road may be modelled by a curve specifying the position in time. The road is the set of places you pass through when you drive. There may be many cars and many speeds, but there is only one road. Local geometry is useful because it isolates the properties of the road which instantaneously affect any car. As it turns out, there are only two such properties for curves in $\mathbb{R}^{3}$ : curvature, quantifying how the curve bends, and torsion, quantifying how it is non-planar. The fundamental theorem of curves states that curvature and torsion classify immersed curves up to rotation and translation.

## 3. Immersed curves, arc-length

$\S 4$. The image of the curve $\alpha(t):=(0,0,0)$ is the single point $(0,0,0)$, not a smooth curve. Another example: let $r \geq 1$ and define

$$
\alpha(t):=\alpha(t)= \begin{cases}\left((-t)^{r+1}, 0,1\right) & t \leq 0 \\ \left(0, t^{r+1}, 1\right) & t \geq 0\end{cases}
$$

This curve is $C^{r}$ and its image is the union of the positive $x$-axis and the positive $y$-axis, shifted up by 1 to avoid obscuring the axes in the figure. The image cannot be regarded as smooth near $\alpha(0)=(0,0,0)$ because there is a corner there. Although the first function is $C^{\infty}$ and the second is $C^{r}$, in neither case is the image a smooth curve. However, supposing $\alpha^{\prime}(a) \neq 0$, the following theorem asserts coordinates which (locally) transform the image of into a line. Under change of coordinates on the codomain, all immersed curves are locally line segments.

$\alpha(t)=(1,1,0)$

$\alpha(t)= \begin{cases}\left((-t)^{r+1}, 0,1\right) & t \leq 0, \\ \left(0, t^{r+1}, 1\right) & t \geq 0 .\end{cases}$


Left and center: the image of a smooth curve may not be a smooth and one dimensional. Right: the image is smooth near $t=a$ such that $\alpha^{\prime}(a) \neq 0$, in the sense of becoming a line segment after a smooth change of coordinates.
§5. Theorem (Local normal form for immersed curves). Let $\alpha$ be a $C^{r}$ curve in $\mathbb{R}^{3}, r \geq 1$, and suppose $\alpha^{\prime}(a) \neq 0$. Then there is a $C^{r}$ diffeomorphism $\varphi: U \subseteq \mathbb{R}^{3} \rightarrow V \subseteq \mathbb{R}^{3}$, and a $\delta>0$, such that $\alpha(t) \in U$ for $a-\delta<t<a+\delta$ and $\varphi \circ \alpha(t)=(t-a, 0,0)$.

Proof (sketch). It is simple enough to use the functional form of the curve to map a line segment on an axis to the curve, eg $\psi(u, v, w)=\left(\alpha_{1}(u), \alpha_{2}(v), \alpha_{3}(w)\right)$, where $\alpha_{1}=\alpha(u+a)$, and $\alpha_{2}(0)=\alpha_{3}(0)=0$. The inverse function theorem can then provide a local inverse $\varphi$ which necessarily maps the curve back to the line segment.
§6. Definition of arc-length: Let $\alpha: I \rightarrow \mathbb{R}^{3}$ and let $a<b \in I$ be such that $[a, b] \subseteq I$. Choose a partition of $[a, b]$ ie choose $t_{i} \in(a, b), i=1, \ldots, n$, such that $a=t_{0}<t_{1}<t_{2}<\ldots<t_{n}=b$. The length of the polygonal path of line segments obtained by joining in sequence the points $\alpha\left(t_{i}\right)$, is $\sum_{i=1}^{n}\left|\alpha\left(t_{i}\right)-\alpha\left(t_{i-1}\right)\right|$. Intuatively, as the $t_{i}$ become close that polygonal path becomes close to the image of $\alpha$. By definition, the length of $\alpha$ between $t=a$ to $t=b$ is length $(\alpha ; a, b)=\lim \sum_{i=1}^{n}\left|\alpha\left(t_{i}\right)-\alpha\left(t_{i-1}\right)\right|$, if that limit exists as the partition $t_{i}$ becomes finer, in the same sense as definition of the Riemann integral. The curve $\alpha$ is called rectifiable if the limit exits for all choices of $a, b$.
$\S 7$. Computation of arc-length: The definition of arc-length is not predicated on differentiability of the curve, but its most common computational formula is.
§8. Proposition (Arc-length formula). If $\alpha: I \rightarrow \mathbb{R}^{3}$ is $C^{1}$, then $\alpha$ is rectifiable and, for all $a, b \in I, a<b$, length $(\alpha ; a, b)=\int_{a}^{b}\left|\alpha^{\prime}(t)\right| d t$.
Proof (sketch). The formula is the limit of the approximation

$$
\sum_{i=1}^{n}\left|\alpha\left(t_{i}\right)-\alpha\left(t_{i-1}\right)\right|=\sum_{i=1}^{n}\left|\frac{\alpha\left(t_{i}\right)-\alpha\left(t_{i-1}\right)}{t_{i}-t_{i-1}}\right|\left(t_{i}-t_{i-1}\right) \approx \sum_{i=1}^{n}\left|\alpha^{\prime}\left(t_{i}^{*}\right)\right|\left(t_{i}-t_{i-1}\right)
$$

where $t_{i}^{*} \in\left[t_{i-1}, t_{i}\right]$, the last term being a Riemann sum for the stated integral.
§9. Reparameterizations: Let $\alpha: I \rightarrow \mathbb{R}^{3}$. A $C^{r}$ reparameterization of $\alpha(t)$ is a $C^{r}$ diffeomorphism $t=t(u), t \in I$. The reparameterized curve is $\alpha(u)=\alpha(t(u))$. Because of the correspondence between $t$ and $u$, the image set, and so the local differential geometry, are the same for a curve and a reparameterization.
§10. Arc length parametrizations: There are preferred, geometrically defined parametrizations, in which the many formulas become much simpler.
$\S 11$. Definition. A curve $\alpha: I \rightarrow \mathbb{R}^{3}$ is parametrized by arc-length if, for all $a, b \in I, a<b, \operatorname{length}(\alpha ; a, b)=a-b$.
$\S 12$. For an arc-length parametrized curve, the parameter, usually denoted " $s$ ", coincides with length along the curve, as though points in the image of the curve are labeled by a ruler bent along the curve. An arc-length reparametrization is one for which the reparametrized curve is parametrized by arc-length. There is a simple test for arc-length parametrization, and, at least in theory, a curve can always be reparametrized by arc-length.
$\S 13$. Proposition. Let $\alpha$ be a $C^{r}$ immersed curve defined on an open interval, $r \geq 1$.

1. $\alpha$ is parametrized by arc-length if and only if it is unit speed ie $\left|\alpha^{\prime}(t)\right|=1$ for all $t$.
2. $s(t)=\int_{a}^{t}\left|\alpha^{\prime}(t)\right| d t$ is a $C^{r}$ diffeomorphism to an open interval, the inverse of which is a $C^{r}$ arc-length reparametrization of $\alpha$.
Proof (sketch). For (1), differentiation of the arc-length formula obtains $\left|\alpha^{\prime}(t)\right|=1$, and the converse is a direct integration. For (2), $s(t)$ is strictly increasing because $s^{\prime}(t)=\left|\alpha^{\prime}(t)\right|>0$, and so $s$ is invertible. The resulting reparameterization is unit speed by the chain rule.
$\S 14$. Arc-length reparametrization is a procedure: find $s$ as a function of $t$, solve that for $t$ as a function of $s$, and then substitute $t=t(s)$ into the function that defined the curve.

## 15. The fundamental theorem of curves

$\S 16$. The best-fit circle to an arc-parametrized curve $\alpha(s)$ at $s=a$ has radius $1 /\left|\alpha^{\prime \prime}(a)\right|$ and center $\alpha(a)+$ $\alpha^{\prime \prime}(a) /\left|\alpha^{\prime \prime}(a)\right|^{2}$. This circle lies in the best fit plane to $\alpha$ at $s=a$, and that plane has normal $\alpha^{\prime}(a) \times \alpha^{\prime \prime}(a)$. See Remarks 40, 42, and Thm. 44.

$\left.T=\alpha^{\prime}(a), N=T^{\prime} /\left|T^{\prime}\right|, B=T \times N\right)$
$\kappa=\left|T^{\prime}\right|, \rho=1 / \kappa, \tau=-N \cdot B^{\prime}$

The best fit circle has radius $\rho$ and center along $N$.
The best fit plane has normal $B$.
Plane curve $\Leftrightarrow \tau=0$.
Serret-Frenet formulas: $T^{\prime}=\kappa N, N^{\prime}=-\kappa T+\tau B, B^{\prime}=\tau N$
Fundamental theorem of curves: $\kappa, \tau \Leftrightarrow \alpha$ up to $S E(3)$
$\boldsymbol{v}=v T, \boldsymbol{a}=\frac{d v}{d t} T+v^{2} \kappa N$,
$T=\frac{\boldsymbol{v}}{v}, B=\frac{\boldsymbol{v} \times \boldsymbol{a}}{|\boldsymbol{v} \times \boldsymbol{a}|}, N=B \times T, \kappa=\frac{|\boldsymbol{v} \times \boldsymbol{a}|}{v^{3}}, \tau=\frac{(\boldsymbol{v} \times \boldsymbol{a}) \cdot \boldsymbol{a}^{\prime}}{|\boldsymbol{v} \times \boldsymbol{a}|^{2}}$
§17. Definition (Serret-Frenet data). Let $\alpha:(a, b) \rightarrow \mathbb{R}^{3}$ be a unit speed curve defined at $t=a$.

1. The unit tangent vector is $T(a)=\alpha^{\prime}(a)$; the unit normal vector is $N(a)=\alpha^{\prime \prime}(a) /\left|\alpha^{\prime \prime}(a)\right|$; the binormal vector is $B(a)=T(a) \times N(a)$.
2. The curvature is $\kappa(a)=\left|\alpha^{\prime \prime}(a)\right|$; the radius of curvature is $\rho(a)=1 / \kappa(a)$.
3. The torsion is $\tau(a)=-N(a) \cdot B^{\prime}(a)$.
4. The Serret-Frenet frame is the tuple $(T(a), N(a), B(a))$. The Serret-Frenet data is the tuple $(\kappa, \tau, T, N, B)$.

For the full Serret-Frenet data, the curve $\alpha$ must be three times differentiable and $\alpha^{\prime \prime}(a)$ must be nonzero.
§18. A best fit circle or plane to a curve can be determined by matching derivatives, as in the Taylor approximations, or from three points on the curve which are made to coalesce. In therms of the Serret-Frenet data, the radius of the circle that best fits the curve at at $t=a$ is $\rho(a)$ and its center at $\alpha(a)+\rho(a) N(a) . B(a)$ is a unit normal to the plane of best fit, and the torsion is zero if that normal is constant. $T$ and $N$ are orthogonal, because $\alpha^{\prime} \cdot \alpha^{\prime \prime}=0$ follows from differentiating $\alpha^{\prime} \cdot \alpha^{\prime}=1$, and $T=\alpha^{\prime}$ and $N=\alpha^{\prime \prime} /\left|\alpha^{\prime \prime}\right|$. Thus, for each $s,(T(s), N(s), B(s))$ is a right-handed orthonormal basis.
§19. If a curve is rotated and translated, then its Serret-Frenet frame is also rotated, because that data is related to best fitting circles, and circles are sent to congruent circles under rotations and translations.
$\S 20$. Theorem (Invariance of the Serret-Frenet data). If $(A, a) \in \mathrm{SE}(3)$ and $\alpha(s)$ is a unit speed curve with SerretFrenet data $(\kappa, \tau, T, N, B)$, then $A \alpha(s)+a$ is unit speed and has Serret-Frenet data ( $A T, A N, A B, \kappa, \tau)$.
Proof (sketch). This is a direct calculation from the formula defining the Serret-Frenet data and from the rotational invariance of the dot and cross product. For example, $\tilde{\alpha}:=A \alpha(s)+a$ is unit speed because $\tilde{\alpha}^{\prime}=A \tilde{\alpha}^{\prime}$ and

$$
\left|\tilde{\alpha}^{\prime}\right|^{2}=\tilde{\alpha}^{\prime} \cdot \tilde{\alpha}^{\prime}=A \alpha^{\prime} \cdot A \alpha^{\prime}=\left|\alpha^{\prime}\right|^{2}=1
$$

$\S 21$. Any vector can be written as a unique sum of vectors in an orthonormal frame. Thus, the derivatives of the Serret-Frenet frame can be decomposed relative to the frame itself, with coefficients that turn out to be the curvature and torsion.
§22. Theorem (Serret-Frenet equations). If $\alpha$ is a $C^{3}$ unit speed curve and $\alpha^{\prime \prime}(t) \neq 0$ for all $t$, then

$$
\frac{d T}{d s}=\kappa N, \quad \frac{d N}{d s}=-\kappa T+\tau B, \quad \frac{d B}{d s}=-\tau N
$$

Proof (sketch). The expansion of aany vector $v$ in terms of $T, N$, and $B$ is $v=(T \cdot v) T+(N \cdot v) N+(B \cdot v) B$. Putting $v=d T / d s$ obtains

$$
\frac{d T}{d s}=\left(T \cdot \frac{d T}{d s}\right) T+\left(N \cdot \frac{d T}{d s}\right) N+\left(B \cdot \frac{d T}{d s}\right) B
$$

and the inner products are easily calculated in terms of $\kappa$ and $\tau$. The formulas for $d N / d s$ and $d B / d s$ are similar.
$\S 23$. Both $\kappa(a)$ and $\tau(a)$ are rates of change of angles with respect to arc-length: $\kappa$ is the the angle of the tangent vector to the normal, and $\tau$ the rate of change of the of the angle between the binormal and the normal (Remark 54). The normal is used as a reference because, by the Serret-Frenet equations, the derivatives of the tangent and the binormal are in the direction of the normal. In these statements, the reference vector is fixed, the angle computed, and only after can the derivative of the angle be computed. The reference vector cannot be allowed to change eg $T$ and $N$ are orthogonal and the derivative of the angle between them is zero. The rate of change of the angles is computed using the dot product and implicit differentiation. It is the derivative $T(s) \cdot N(a)$ ( $a$ fixed) that is involved, not the derivative of $T(s) \cdot N(s)=1$.
$\S 24$. Besides being important computationally, the Serret-Frenet equations are important because, given $\kappa$ and $\tau$, they are differential equations for $T, N$, and $B$, thus enabling an application of the powerful existence on uniqueness theory of ordinary differential equations.
$\S 25$. Theorem (Fundamental theorem of curves). Let $r \geq 3$, and let $\kappa_{0}>0$ and $\tau_{0}$ be $C^{r-2}$ and $C^{r-3}$ functions on and open interval $I$, respectively. Then there is a $C^{r}$ immersed curve $\alpha: I \rightarrow \mathbb{R}^{3}$ with curvature $\kappa=\kappa_{0}$ and torsion $\tau=\tau_{0}$. If $\alpha_{1}$ and $\alpha_{2}$ are two such curves, then there is an element $(A, a) \in \mathrm{SE}(3)$ such that $\alpha_{2}(s)=A \alpha_{1}(s)+a$.
Proof (sketch). The initial value problem of the Serret-Frenet formulas, with initial data the standard coordinate orthonormal frame, has a unique solution. The $T, N$, and $B$ so obtained are an orthonormal frame because the relevant inner products satisfy a (different) initial value problem which has a (unique) constant solution. $\alpha^{\prime}=T$, so the curve $\alpha$ is constructed from $T$ by integration, and has the given curvature and torsion because those involve higher than first derivatives of $\alpha$, so at least the first derivatives of the constructed $T, N$, and $B$, and those derivatives are determined by their defining differential equation. The uniqueness of the curve up to $\mathrm{SE}(3)$ follows from the uniqness of initial value problems and the invariance of the Serret-Frenet data.
§26. Special curves: The fundamental theorem of curves determines the curve from the curvature and torsion, thus classifying curves up to rotation and translation as pairs of functions. Because of this classification, the question can be put as to whether some of the simpler curves can be recognized from the pair. In summary, a line is characterized by zero curvature, a plane curve by zero torsion, and a circle by constant curvature. A general helix, moving along an axis at a constant pitch, is characterized by a constant ratio of curvature and torsion. A curve on a sphere has a more-or-less complicated functional relative between curvature and torsion.

## §27. Theorem.

1. A line is an immersed curve in $\mathbb{R}^{3}$ which can be reparametrized to $\alpha(t)=a+b t$, where $a$ and $b$ are constant vectors, $b \neq 0$. For $C^{2}$ immersed curves in $\mathbb{R}^{3}$, the follow statements are equivalent: (A) $\alpha$ is a line; ( $B$ ) $T$ is constant; (C) $\kappa=0$.
2. An immersed curve $\alpha(t)$ in $\mathbb{R}^{3}$ is planar if there is a nonzero constant vector $a$ such that $\alpha(t) \cdot a$ is constant. For $C^{3}$ immersed curves $\alpha$ in $\mathbb{R}^{3}$ with $\kappa$ never zero, the following statements are equivalent: (A) $\alpha$ is planar; (B) $B$ is constant; (C) $\tau=0$.
3. An immersed curve $\alpha(t)$ in $\mathbb{R}^{3}$ is a sphere curve (with center $a$ ) if there is a constant nonzero vector a such that $|\alpha(t)-a|$ is constant. For $C^{4}$ immersed curves $\alpha$ in $\mathbb{R}^{3}$ with $\kappa$ never zero and $\tau$ never zero, the following statements are equivalent: (A) $\alpha$ is a sphere curve; $(B) \alpha+\rho N+\left(\rho^{\prime} \sigma / v\right) B$ is constant; $(C)\left(\sigma \rho^{\prime} / v\right)^{\prime} / v+\rho \tau=0$.
4. A circle is a planar sphere curve. For $C^{3}$ immersed curves $\alpha$ in $\mathbb{R}^{3}$ with $\kappa$ never zero, the following statements are equivalent: ( $A$ ) $\alpha$ is a circle; $(B) \alpha+\rho N$ is constant and $B$ is constant; (C) $\kappa$ is constant and $\tau=0$.
5. A helix (with pitch $|a|$ and axial vector $a$ ) is an immersed curve in $\mathbb{R}^{3}$ which can be reparametrized to $\alpha(t)=\bar{\alpha}(t)+t a$, where $a \neq 0$ is constant and $\bar{\alpha}$ is unit speed and in a plane with normal $a$. For $C^{3}$ immersed curves $\alpha$ in $\mathbb{R}^{3}$ with $\kappa$ never zero and $\tau$ never zero, the following statements are equivalent: (A) $\alpha$ is a helix; (B) there is a nonzero constant vector $a$ such that $T \cdot a$ is a nonzero constant; ( $C$ ) $\tau / \kappa$ is a nonzero constant.
6. A circular helix is an immersed curve in $\mathbb{R}^{3}$ which can be reparametrized to $\alpha(t)=\bar{\alpha}(t)+a t$ where where $a \neq 0$ is constant and $\bar{\alpha}(t)$ is a circle in a plane with normal $a$. For $C^{3}$ immersed curves $\alpha$ in $\mathbb{R}^{3}$ with $\kappa$ never zero and $\tau$ never zero, the following statements are equivalent: (A) $\alpha$ is a circular helix; (B) there is a nonzero vector $a$ such that $T \cdot a$ is a nonzero constant, and $\left(\alpha+\rho v^{2} N\right) \times a$ is constant; ( $C$ ) $\kappa$ and $\tau$ are nonzero constants.
§28. Calculation of Serret-Frenet data: For a curve $\alpha(t)$, not necessarily unit speed, the velocity is $\boldsymbol{v}:=\alpha^{\prime}$; the speed is $v:=|\boldsymbol{v}|$; the acceleration is $\boldsymbol{a}:=\alpha^{\prime \prime}$; the scalar acceleration is $a:=|\boldsymbol{a}|$, and the Serret-Frenet data of the corresponding arc-length parameterized curve can be calculated from these, without explicitly calculating the arc-length. Conversely, the velocity $\boldsymbol{v}$ and acceleration $\boldsymbol{a}$ can be written in terms of the Serret-Frenet frame.
$\S 29$. Theorem. For $C^{3}$ immersed curve $\alpha(t)$ (not necessarily unit speed) such that $\alpha^{\prime \prime}(t) \neq 0$ for all $t$,

$$
\boldsymbol{v}=v T, \quad \boldsymbol{a}=\frac{d v}{d t} T+v^{2} \kappa N, \quad T=\frac{\boldsymbol{v}}{v}, \quad B=\frac{\boldsymbol{v} \times \boldsymbol{a}}{|\boldsymbol{v} \times \boldsymbol{a}|}, \quad N=B \times T, \quad \kappa=\frac{|\boldsymbol{v} \times \boldsymbol{a}|}{v^{3}}, \quad \tau=\frac{(\boldsymbol{v} \times \boldsymbol{a}) \cdot \boldsymbol{a}^{\prime}}{|\boldsymbol{v} \times \boldsymbol{a}|^{2}},
$$

where $(\kappa, \tau, T, N, B)$ is the Serret-Frenet data of any arc-length parameterization.
Proof (sketch). The arc-length parametrized curve $\alpha(s)$ is related to the original curve $\alpha(t)$ by the equation $\alpha(t)=\alpha(s(t))$. The chain rule calculates the derivatives through this substitution eg $\boldsymbol{v}=d \alpha / d t=d \alpha / d s d s / d t=$ $\boldsymbol{v}(t) T(s(t))$, and the formulas are obtained by further calculating the second and third derivatives of $\alpha(t)$ using the Serret-Frenet formulas.
$\S 30$. The formulas provide insight into three dimensional motion. The curvature is zero where $\boldsymbol{v}$ and $\boldsymbol{a}$ are parallel, and the torsion is arising from the component along $B$ of the third derivative of the curve. The velocity $\boldsymbol{v}$ and acceleration $\boldsymbol{a}$ are in the best-fit plane, because they are both orthogonal to $B$. From the second formula, the tangent component of the acceleration is the rate of change of speed, and its normal component depends only on the speed and the curvature. This last aspect is a generalization of the common formula that the centripetal acceleration of an object in a circular motion is $v^{2} / R$. The scalar acceleration is not the (unsigned) rate of change of speed unless either the curvature is zero or the speed is zero, since $a=|\boldsymbol{a}|=\sqrt{(d v / d t)^{2}+\left(v^{2} \kappa\right)^{2}}$.

## 31. Epilogue

$\S 32$. The study of curves is the beginning of differential geometry, because curves, being dimension 1 , are the simplest nontrivial objects. The result is satisfyingly complete. Essentially, all curves are classified, up to congruence, by two functions - curvature and torsion. The correspondence is constructive: the curve corresponding to given curvature $\kappa$ and torsion $\tau$ can be found by solving the initial value problem

$$
\frac{d \boldsymbol{x}}{d s}=T, \quad \frac{d T}{d s}=\kappa N, \quad \frac{d N}{d s}=-\kappa T+\tau B, \quad \frac{d B}{d s}=\tau N, \quad \boldsymbol{x}(0)=0, T(0)=\boldsymbol{i}, N(0)=\boldsymbol{j}, B(0)=\boldsymbol{k}
$$

and such can be numerically solved to practically any precision. The development identifies quantities-the SerretFrenet data- that are specific to the geometry of the curve, as opposed to being of the way that the curve is moved along. With those identified, motion along three dimensional curves can be separated into a geometric and dynamic parts, and motions of objects in three space are better understood.
§33. From the modest study of curves emerges far-reaching themes of differential geometry. Immersed curves are locally line segments - and generally in Differential Geometry, the stance is that smooth at a point means locally diffeomorphic to linear, and otherwise the point is a singularity. An arc-length parametrization can be thought of as an assignment of a number to each point of a curve, at least locally (because of the possibility that the curve might self intersect). This is the most elementary example of a special, geometrically motivated coordinate system, in
which formulas greatly simplify. The most significant progress was made through the use of two powerful theorems of analysis: the inverse function theorem and the existence and uniqueness of solutions of odes.
$\S 34$. Another emergent theme is the use of geometry to guide and motivate definitions and analysis. For example, to find the curvature and torsion, circles were fit to curves in geometrically meaningful ways: by best fit according to a geometrically meaningful distance, or by passing circles through coalescing points. Formally, none of this is required: simply define the Serret-Frenet data according to the formulas in Definition 17, and the development will be shorter. There is no real need to derive the arc-length formula from the more basic Riemann sums; the formula itself would adequately serve as the definition. The cost of abandoning geometric motivations is a loss of vision, and that can lead to loss of simplicity and elegance. Differential geometry brings vision to analysis. That is present even in the simplest study.

## 35. Technicalities

$\S 36$. Theorem (Local normal form for immersed curves). Let $\alpha$ be a $C^{r}$ curve in $\mathbb{R}^{3}, r \geq 1$, and suppose $\alpha^{\prime}(a) \neq 0$. Then there is a $C^{r}$ diffeomorphism $\varphi: U \subseteq \mathbb{R}^{3} \rightarrow V \subseteq \mathbb{R}^{3}$, and a $\delta>0$, such that $\alpha(t) \in U$ for $a-\delta<t<a+\delta$ and $\varphi \circ \alpha(t)=(t-a, 0,0)$.

Proof. Define the map $\psi$ by

$$
\psi(u, v, w)=\alpha(u+a)+\alpha_{2}(v)+\alpha_{3}(w)
$$

where $\alpha_{2}$ and $\alpha_{3}$ are any $C^{k}$ curves of $(a, b)$ such that $\alpha_{2}(0)=\alpha_{3}(0)=0$ and $\alpha^{\prime}(a), \alpha_{2}^{\prime}(0), \alpha_{3}^{\prime}(0)$ is a basis of $\mathbb{R}^{3}$. For example, one can complete the vector $\alpha^{\prime}(a)$ to an orthonormal basis $\alpha^{\prime}\left(a, e_{2}, e_{3}\right.$ and choose $\alpha_{2}(t)=t e_{2}$ and $\alpha_{3}(t)=t e_{3}$. Then $\psi(0,0,0)=\alpha(a)$ and the derivative of $\psi$ at $(0,0,0)$ is the $3 \times 3$ matrix $\left[\alpha^{\prime}(a), \alpha_{2}^{\prime}(0), \alpha_{3}^{\prime}(0)\right]$, the rows of which are linearly independent. By the inverse function theorem, $\psi$ has a $C^{k}$ local inverse $\varphi: U \ni \alpha(a) \rightarrow$ $V \ni(0,0,0)$. Since $V$ is open, there is a $\delta>0$ such that $(u, 0,0) \in V$ whenever $u \in(-\delta, \delta)$, and for such a $u$, $(u, 0,0)=\varphi(\psi(u, 0,0))=\varphi(\alpha(u+a))$. Substituting $t=u-a$ obtains $\varphi(\alpha(t))=(t+a, 0,0)$ for $t \in(a-\delta, a+\delta)$.
§37. Proposition (Arc-length formula). If $\alpha: I \rightarrow \mathbb{R}^{3}$ is $C^{1}$, then $\alpha$ is rectifiable and, for all $a, b \in I$, $a<b$, length $(\alpha ; a, b)=\int_{a}^{b}\left|\alpha^{\prime}(t)\right| d t$.
Proof. By the definition of derivative, $\alpha(t+u)-\alpha(t)=\alpha^{\prime}(t) u+R(t, u)$ where $\lim _{u \rightarrow 0} R(t, u) / u=0$.

$$
|\alpha(t+u)-\alpha(t)|=\left|\alpha^{\prime}(t)+\frac{R(t, u)}{u}\right| u=\left|\alpha^{\prime}(u)\right|+e(t, u), \quad e(t, u):=\left|\alpha^{\prime}(t)+\frac{R(t, u)}{u}\right|-\left|\alpha^{\prime}(t)\right|
$$

Since $\left|\alpha^{\prime}\right|$ is continuous, it is Riemann integrable on $[a, b]$, and also $\lim _{u \rightarrow 0} e(t, u)=0$, so by Lem. 86,

$$
\begin{aligned}
\lim \sum_{i}\left|\alpha\left(t_{i}\right)-\alpha\left(t_{i-1}\right)\right| & =\lim \sum_{i}\left(\left|\alpha^{\prime}\left(t_{i-1}\right)\right| \Delta t_{i}+e\left(t_{i}, \Delta t_{i}\right)\right) \\
& =\lim \sum_{i}\left|\alpha^{\prime}\left(t_{i-1}\right)\right| \Delta t_{i}+\lim \sum_{i} e\left(t_{i}, \Delta t_{i}\right)=\int_{p}^{q}\left|\alpha^{\prime}(t)\right| d t
\end{aligned}
$$

§38. Lemma. A strictly increasing $C^{r}$ function, $r \geq 1$, that is defined on an open interval and has positive derivative, is a $C^{r}$ diffeomorphism between open intervals.
Proof. Let $f:(a, b) \rightarrow \mathbb{R}$ be continuous with positive derivative, and let $c:=\inf f(a, b)$ and $d:=\sup f(a, b)$. If $y \in(c, d)$ then $y>c[y<d]$ and there is a $y^{-} \in f(a, b), y^{-}=f\left(x^{-}\right)\left[y^{+} \in f(a, b), y^{+}=f\left(x^{+}\right)\right]$such that $c<y<y^{-}$ $\left[y<y^{+}<d\right]$. Then $y^{-}<y<y^{+}$and by the intermediate value theorem there is an $x$ such that $x^{-}<x<x^{+}$and $f(x)=y$, so $f$ is onto $(c, d)$. If $f\left(x_{1}\right) \neq f\left(x_{2}\right)$ then either $x_{1}<x_{2}$ or $x_{1}>x_{2}$ and $f\left(x_{1}\right)<f\left(x_{2}\right)$ or $f\left(x_{1}\right)>f\left(x_{2}\right)$, respectively, and in either case $f\left(x_{1}\right) \neq f\left(x_{2}\right)$, so $f$ is bijective. The inverse of $f$ is $C^{r}$ by the inverse function theorem.
§39. Proposition. Let $\alpha$ be a $C^{r}$ immersed curve defined on an open interval, $r \geq 1$.

1. $\alpha$ is parametrized by arc-length if and only if it is unit speed ie $\left|\alpha^{\prime}(t)\right|=1$ for all $t$.
2. $s(t)=\int_{a}^{t}\left|\alpha^{\prime}(t)\right| d t$ is a $C^{r}$ diffeomorphism to an open interval, the inverse of which is a $C^{r}$ arc-length reparametrization of $\alpha$.

Proof. For (1), if $\alpha$ is parametrized by arc-length and $a \in \operatorname{domain}(\alpha)$, then

$$
1=\left.\frac{d}{d s}\right|_{s=p}(s-a)=\left.\frac{d}{d s}\right|_{s=p} \int_{a}^{s}\left|\alpha^{\prime}(a)\right| d t=\left|\alpha^{\prime}(a)\right|
$$

and conversely, supposing $\left|\alpha^{\prime}(t)\right|=1$ for all $t$,

$$
\operatorname{length}(\alpha ; a, b)=\int_{a}^{b}\left|\alpha^{\prime}(t)\right| d t=\int_{a}^{b} 1 d t=a-b
$$

For (2), $s(t)$ is $C^{1}$, and $s^{\prime}(t)=\left|\alpha^{\prime}(t)\right|$ by the fundamental theorem of calculus, so $s(t)$ is $C^{r}$ since $\left|\alpha^{\prime}(t)\right|$ is $C^{r-1}$. By Lem. $38 s(t)$ is a $C^{r}$ diffeomorphism to an open interval. The inverse $t(s)$ can be used to reparametrize $\alpha(t)$, and that is an arc-length reparametrization because, by the chain rule

$$
s^{\prime}(t)=\left|\frac{d \alpha}{d t}\right|=\left|\frac{d \alpha}{d s} \frac{d s}{d t}\right|=\left|\frac{d \alpha}{d s}\right| s^{\prime}(t)
$$

and $s^{\prime}(t) \neq 0$ implies $|d \alpha / d s|=1$.


Left: parametrization of a circle in general orientation in $\mathbb{R}^{3}$. Center: fitting that circle using the error defined as the distance between equal-arc points on the circle and the curve. Right: fitting by using three coalescing points.
$\S 40$. Remark (Fitting a circle to a given curve). Let $\hat{a}$ and $\hat{b}$ be orthogonal unit vectors, $\rho>0$, and $c$ be a point in $\mathbb{R}^{3}$. The (unit speed) curve

$$
\beta(s)=c+\rho \cos (s / \rho) \hat{a}+\rho \sin (s / \rho) \hat{b}
$$

is a circle of radius $\rho$, centered at $c$, in the plane which is parallel to $\hat{a}$ and $\hat{b}$. Because the interest is about the set $\{\alpha(s) \mid s \in I\}$, its is necessary to fit that set to the set $\{\beta(s) \mid s \in I\}$, near $\alpha(a)$. One way to do this is to measure the error between the two curves as the square of the distance $e(s)$ between points of the curve at equal-arc distances from $\alpha(a)$, ie

$$
e(s)=|\boldsymbol{e}(s)|^{2}, \quad \boldsymbol{e}(s)=\alpha(s)-\beta(s-a)
$$

and set as many as possible derivatives of $e(s)$ to zero, because that corresponds to the smallest error for small (Figure 2, left and center). Since

$$
e=\boldsymbol{e} \cdot \boldsymbol{e}, \quad e^{\prime}=2 \boldsymbol{e} \cdot \boldsymbol{e}^{\prime}, \quad e^{\prime \prime}=2 \boldsymbol{e} \cdot \boldsymbol{e}^{\prime \prime}+2 \boldsymbol{e}^{\prime} \cdot \boldsymbol{e}^{\prime}, \quad e^{\prime \prime \prime}=2 \boldsymbol{e} \cdot \boldsymbol{e}^{\prime \prime \prime}+6 \boldsymbol{e}^{\prime} \cdot \boldsymbol{e}^{\prime \prime}, \quad e^{(4)}=2 \boldsymbol{e} \cdot \boldsymbol{e}^{(4)}+8 \boldsymbol{e}^{\prime} \cdot \boldsymbol{e}^{\prime \prime \prime}+6 \boldsymbol{e}^{\prime \prime} \cdot \boldsymbol{e}^{\prime \prime}
$$

this is the same as setting $\boldsymbol{e}(0)=\boldsymbol{e}^{\prime}(a)=\boldsymbol{e}^{\prime \prime}(a)=0\left(e(a)=0\right.$ is the same as $\boldsymbol{e}(a)=0$, and then $e^{\prime}(a)=0$ automatically, from which $e^{\prime \prime}=0$ is the same as $\boldsymbol{e}^{\prime}(a)=0$, and so on). So for the best fit of the two sets amounts to derivative matching of curves in the arc-length parametrization:

$$
\begin{equation*}
\boldsymbol{e}(a)=0 \Leftrightarrow \alpha(a)=c+\rho \hat{a}, \quad \boldsymbol{e}^{\prime}(a)=0 \Leftrightarrow \alpha^{\prime}(a)=\hat{b}, \quad \boldsymbol{e}^{\prime \prime}(a)=0 \Leftrightarrow \alpha^{\prime \prime}(a)=(-1 / \rho) \hat{a} \tag{41}
\end{equation*}
$$

Equations (41) are consistent with $a \cdot b=0$, because $\alpha^{\prime} \cdot \alpha^{\prime \prime}=0$, since ' is the derivative with respect to arc-length. From the third of (41), $\rho=1 /\left|\alpha^{\prime \prime}(a)\right|, \kappa=\left|\alpha^{\prime \prime}(a)\right|$, and $\hat{a}=-\alpha^{\prime \prime}(a) /\left|\alpha^{\prime \prime}(a)\right|$. The second of (41) implies $\hat{b}=\alpha^{\prime}(0)$, and the first implies the center of the circle is located at $c=\alpha(a)-\rho \hat{a}=\alpha(a)+\alpha^{\prime \prime}(a) /\left|\alpha^{\prime \prime}(a)\right|^{2}$.
$\S 42$. Remark (Fitting a plane to a given curve). Let $a$ be a nonzero vector and $c$ be a point in $\mathbb{R}^{3}$. The plane perpendicular to $a$ through the point $c$ is $a \cdot(x-c)=0$. To best-fit this plane to a given immersed curve $\alpha(t)$, measure the error between the curve and the plane as the signed distance $e(t)$ from points on the curve to the plane, i.e., $e(t)=(\alpha(t)-c) \cdot a$, and set as many as possible derivatives of $e(s)$ to zero, because that corresponds to the smallest error for small $t$ :

$$
e(a)=(\alpha(a)-c) \cdot a=0, \quad e^{\prime}(a)=\alpha^{\prime}(a) \cdot a=0, \quad e^{\prime \prime}(a)=\alpha^{\prime \prime}(a) \cdot a=0
$$

The point $c$ is ambiguous up to addition of vectors orthogonal to $a$, so one can take $c=\alpha(a)$. The next two equations imply $a$ is orthogonal to $\alpha^{\prime}(a)$ and $\alpha^{\prime \prime}(a)$, and one can take $a=\alpha^{\prime}(a) \times \alpha^{\prime \prime}(a)$.
§43. Lemma. There is a unique circle containing $a_{1}, a_{2} \in \mathbb{R}^{3}$ and the origin if these three points are not collinear. The center of that circle is at $\gamma a_{1}+\delta a_{2}$, where

$$
\gamma=\frac{1}{2 D}\left|a_{2}\right|^{2}\left(\left|a_{1}\right|^{2}-a_{1} \cdot a_{2}\right), \quad \delta=\frac{1}{2 D}\left|a_{1}\right|^{2}\left(\left|a_{2}\right|^{2}-a_{1} \cdot a_{2}\right), \quad D=\left|a_{1} \times a_{2}\right|^{2}
$$

Proof. Let the center of the circle be at $c$. As the circle and its center lie in a the same plane, $c=\gamma a_{1}+\delta a_{2}$ for some $\gamma, \delta$. Since $a_{1}, 0$, and $a_{2}$ are all on the same circle, $\left|a_{1}-c\right|^{2}=|c|^{2}=\left|a_{2}-c\right|^{2}$, implying

$$
0=\left|a_{1}-c\right|^{2}-|c|^{2}=\left|a_{1}\right|^{2}-2 a_{1} \cdot c+|c|^{2}-|c|^{2}=\left|a_{1}\right|^{2}-2 a_{1} \cdot c
$$

so $a_{1} \cdot c=\left|a_{1}\right|^{2} / 2$, and similarly $a_{2} \cdot c=\left|a_{2}\right|^{2} / 2$. By dotting the equation $c=\gamma a_{1}+\delta a_{2}$ with $a_{1}$ and then $a_{2}$,

$$
\frac{\left|a_{1}\right|^{2}}{2}=a_{1} \cdot c=\gamma\left|a_{1}\right|^{2}+\delta a_{1} \cdot a_{2}, \quad \frac{\left|a_{2}\right|^{2}}{2}=a_{1} \cdot c=\gamma a_{1} \cdot a_{2}+\delta\left|a_{2}\right|^{2}
$$

The result is obtained by solving these linear equations for $\gamma$ and $\delta$.
$\S 44$. Theorem. Let $\alpha$ be a $C^{2}$ unit speed curve and $\alpha^{\prime \prime}(a) \neq 0$.

1. $\lim _{h \rightarrow 0^{+}} c(h)=\alpha(a)+\alpha^{\prime \prime}(a) /\left|\alpha^{\prime \prime}(a)\right|^{2}$, where $c(h)$ is the center of the circle through $\alpha(a-h), \alpha(a)$, and $\alpha(a+h)$.
2. $\lim _{h \rightarrow 0} n(h)=\alpha^{\prime}(a) \times \alpha^{\prime \prime}(a) /\left|\alpha^{\prime}(a) \times \alpha^{\prime \prime}(a)\right|$, where $n(h)$ is the unit normal of the plane passing though $\alpha(p-h), \alpha(a)$, and $\alpha(p+h)$, defined as

$$
n(h):=\frac{a_{1}(h) \times a_{2}(h)}{\left|a_{1}(h) \times a_{2}(h)\right|}, \quad a_{1}(h):=\alpha(p-h)-\alpha(a), \quad a_{2}(h):=\alpha(p+h)-\alpha(a) .
$$

Proof. By translating the curve, assume without loss of generality that $\alpha(a)=0$. Let $a_{1}(h)=\alpha(s-h)$ and $a_{2}(h)=$ $\alpha(s+h)$. Then Lem. 43 provides $\gamma(h), \delta(h), D(h)$, and $c(h)$, and it is required to show that $\lim _{h \rightarrow 0^{+}} c(h)=\frac{\alpha^{\prime \prime}(a)}{\left|\alpha^{\prime \prime}(a)\right|^{2}}$. By Taylor's formula, $\alpha(p+h)=\alpha(a)+\alpha^{\prime}(a) h+\frac{1}{2} \alpha^{\prime \prime}(a) h^{2}+o\left(h^{2}\right)$, so, dropping the explicit dependence of functions on $a\left(\operatorname{eg} \alpha^{\prime}\right.$ without evaluation means $\left.\alpha^{\prime}(a)\right)$,

$$
\begin{aligned}
& a_{2}(h)=\alpha^{\prime} h+\frac{1}{2} \alpha^{\prime \prime} h^{2}+o\left(h^{2}\right), \quad a_{1}(h)=-\alpha^{\prime} h+\frac{1}{2} \alpha^{\prime \prime} h^{2}+o\left(h^{2}\right) \\
& c(h)=\gamma(h) a_{1}(h)+\delta(h) a_{2}(h)=\frac{h^{2}}{2}(\delta(h)+\gamma(h)) \alpha^{\prime \prime}+(\delta(h)-\gamma(h)) h \alpha^{\prime}+\delta(h) o\left(h^{2}\right)+\gamma(h) o\left(h^{2}\right)
\end{aligned}
$$

and it suffices to show

$$
\lim _{h \rightarrow 0^{+}} \frac{h^{2}}{2}(\delta(h)+\gamma(h))=\frac{1}{\left|\alpha^{\prime \prime}\right|^{2}}, \quad \lim _{h \rightarrow 0^{+}}(\delta(h)-\gamma(h)) h=0, \quad \lim _{h \rightarrow 0^{+}} \delta(h) o\left(h^{2}\right)=0, \quad \lim _{h \rightarrow 0^{+}} \delta(h) o\left(h^{2}\right)=0
$$

$\S 45$. Substitute $a_{1}(h)$ and $a_{2}(h)$, remembering $\alpha^{\prime}(a) \cdot \alpha^{\prime}(a)=1$ and $\alpha^{\prime}(a) \cdot \alpha^{\prime \prime}(a)=0$ :

$$
\begin{aligned}
& \left|a_{2}\right|^{2}=\left|\alpha^{\prime} h\right|^{2}+\frac{1}{4}\left|\alpha^{\prime \prime} h^{2}\right|^{2}+o\left(h^{4}\right)+\left(\alpha^{\prime} h\right) \cdot\left(\alpha^{\prime \prime} h^{2}\right)+\left(\alpha^{\prime} h\right) \cdot o\left(h^{2}\right)+\left(\alpha^{\prime \prime} h^{2}\right) \cdot o\left(h^{2}\right)=h^{2}+o\left(h^{3}\right), \\
& \left|a_{1}\right|^{2}=h^{2}+o\left(h^{3}\right), \\
& a_{1} \cdot a_{2}=-h^{2}+o\left(h^{3}\right), \\
& a_{1} \times a_{2}=\alpha^{\prime} \times \alpha^{\prime \prime} h^{3}+o\left(h^{3}\right), \\
& D=\left|a_{1} \times a_{2}\right|^{2}=\left|\alpha^{\prime} \times \alpha^{\prime \prime}\right|^{2} h^{6}+o\left(h^{6}\right)=\left|\alpha^{\prime \prime}\right|^{2} h^{6}+o\left(h^{6}\right), \\
& \gamma=\frac{1}{2 D}\left|a_{2}\right|^{2}\left(\left|a_{1}\right|^{2}-a_{1} \cdot a_{2}\right)=\frac{\left(h^{2}+o\left(h^{3}\right)\right)\left(2 h^{2}+o\left(h^{3}\right)\right)}{2\left(\left|\alpha^{\prime \prime}\right|^{2} h^{6}+o\left(h^{6}\right)\right)}=\frac{h^{4}+o\left(h^{5}\right)}{\left|\alpha^{\prime \prime}\right|^{2} h^{6}+o\left(h^{6}\right)}=\frac{1+o(h)}{\left|\alpha^{\prime \prime}\right|^{2} h^{2}}, \\
& \delta=\frac{1+o(h)}{\left|\alpha^{\prime \prime}\right|^{2} h^{2}}, \\
& \lim _{h \rightarrow 0^{+}} \frac{h^{2}}{2}(\delta(h)+\gamma(h))=\lim _{h \rightarrow 0^{+}} \frac{h^{2}}{2}\left(\frac{1+o(h)}{\left|\alpha^{\prime \prime}\right|^{2} h^{2}}+\frac{1+o(h)}{\left|\alpha^{\prime \prime}\right|^{2} h^{2}}\right)=\frac{1}{\left|\alpha^{\prime \prime}\right|^{2}} \\
& \lim _{h \rightarrow 0^{+}}(\delta-\gamma) h=\lim _{h \rightarrow 0^{+}} \frac{o\left(h^{2}\right)}{\left|\alpha^{\prime \prime} h^{2}\right|}=0 \\
& \lim _{h \rightarrow 0^{+}} \gamma(h) o\left(h^{2}\right)=\lim _{h \rightarrow 0^{+}} \frac{o\left(h^{2}\right)+o\left(h^{3}\right)}{\left|\alpha^{\prime \prime}\right|^{2} h^{2}}=0, \\
& \lim _{h \rightarrow 0^{+}} \delta(h) o\left(h^{2}\right)=\lim _{h \rightarrow 0^{+}} \frac{o\left(h^{2}\right)+o\left(h^{3}\right)}{\left|\alpha^{\prime \prime}\right|^{2} h^{2}}=0 .
\end{aligned}
$$

For the second statement,

$$
\begin{aligned}
\lim _{h \rightarrow 0^{+}} n(h) & =\lim _{h \rightarrow 0^{+}} \frac{a_{1}(h) \times a_{2}(h)}{\left|a_{1}(h) \times a_{2}(h)\right|} \\
& =\lim _{h \rightarrow 0^{+}} \frac{\alpha^{\prime} \times \alpha^{\prime \prime} h^{3}+o\left(h^{3}\right)}{\sqrt{\left|\alpha^{\prime} \times \alpha^{\prime \prime}\right|^{2} h^{6}+o\left(h^{6}\right)}} \\
& =\lim _{h \rightarrow 0^{+}} \frac{\alpha^{\prime} \times \alpha^{\prime \prime}+o\left(h^{0}\right)}{\sqrt{\left|\alpha^{\prime} \times \alpha^{\prime \prime}\right|^{2}+o\left(h^{0}\right)}} \\
& =\frac{\alpha^{\prime}(a) \times \alpha^{\prime \prime}(a)}{\left|\alpha^{\prime}(a) \times \alpha^{\prime \prime}(a)\right|} .
\end{aligned}
$$

$\S 46$. Theorem (Invariance of the Serret-Frenet data). If $(A, a) \in \operatorname{SE}(3)$ and $\alpha(s)$ is a unit speed curve with SerretFrenet data $(\kappa, \tau, T, N, B)$, then $A \alpha(s)+a$ is unit speed and has Serret-Frenet data ( $A T, A N, A B, \kappa, \tau)$.
Proof. Let $\tilde{\alpha}(s)=A \alpha(s)+a$ and let the Serret-Frenet data of $\tilde{\alpha}$ be denoted similarly. Then

$$
\tilde{\alpha}^{\prime}(s) \cdot \tilde{\alpha}^{\prime}(s)=\left(A \alpha^{\prime}(s)\right) \cdot\left(A \alpha^{\prime}(s)\right)=\alpha^{\prime}(s) \cdot \alpha^{\prime}(s)=1
$$

so $\tilde{\alpha}(s)$ is unit speed, and

$$
\begin{aligned}
& \tilde{T}(s)=\tilde{\alpha}^{\prime}(s)=A \alpha^{\prime}(s)=A T(s) \\
& \tilde{\kappa}=\left|\tilde{T}^{\prime}(s)\right|=\left|A T^{\prime}(s)\right|=\left|T^{\prime}(s)\right|=\kappa \\
& \tilde{N}(s)=\tilde{T}^{\prime}(s) / \tilde{\kappa}(s)=A T^{\prime}(s) / \kappa(s)=A N(s) \\
& \tilde{B}(s)=\tilde{T}(s) \times \tilde{N}(s)=(A T(s)) \times(A N(s))=A(T(s) \times N(s))=A B(s) \\
& \tilde{\tau}(s)=-\tilde{N}(s) \cdot \tilde{B}^{\prime}(s)=-(A N(s)) \cdot(A B(s))^{\prime}=-(A N(s)) \cdot\left(A B^{\prime}(s)\right)=-N(s) \cdot B^{\prime}(s)=\tau(s)
\end{aligned}
$$

$\S 47$. Theorem (Serret-Frenet equations). If $\alpha$ is a $C^{3}$ unit speed curve and $\alpha^{\prime \prime}(t) \neq 0$ for all $t$, then

$$
\frac{d T}{d s}=\kappa N, \quad \frac{d N}{d s}=-\kappa T+\tau B, \quad \frac{d B}{d s}=-\tau N
$$

Proof. Let $u_{1}:=T, u_{2}:=N, u_{3}:=B$. By the expansion (63), $u_{i}^{\prime}=\sum_{j=1}^{3} a_{i j} u_{j}$ where $a_{i j}=u_{i}^{\prime} \cdot u_{j}$, and it is required to show that

$$
\left[a_{i j}\right]=\left[\begin{array}{ccc}
0 & \kappa & 0 \\
-\kappa & 0 & \tau \\
0 & -\tau & 0
\end{array}\right]
$$

The matrix $\left[a_{i j}\right]$ is antisymmetric because $0=\left(u_{i} \cdot u_{j}\right)=u_{i}^{\prime} \cdot u_{j}+u_{i} \cdot u_{j}^{\prime}=a_{i j}+a_{j i}$, so the diagonal elements are zero, and

$$
a_{12}=N \cdot T^{\prime}=N \cdot(\kappa N)=\kappa N \cdot N=\kappa, \quad a_{13}=B \cdot T^{\prime}=B \cdot(\kappa N)=0, \quad a_{23}=-a_{32}=-N \cdot B^{\prime}=\tau
$$

§48. Theorem (Fundamental theorem of curves). Let $r \geq 3$, and let $\kappa_{0}>0$ and $\tau_{0}$ be $C^{r-2}$ and $C^{r-3}$ functions on and open interval $I$, respectively. Then there is a $C^{r}$ immersed curve $\alpha: I \rightarrow \mathbb{R}^{3}$ with curvature $\kappa=\kappa_{0}$ and torsion $\tau=\tau_{0}$. If $\alpha_{1}$ and $\alpha_{2}$ are two such curves, then there is an element $(A, a) \in \operatorname{SE}(3)$ such that $\alpha_{2}(s)=A \alpha_{1}(s)+a$.
Proof. The linear differential equation

$$
\begin{equation*}
\frac{d u_{1}}{d s}=\kappa_{0}(s) u_{2}, \quad \frac{d u_{2}}{d s}=-\kappa_{0}(s) u_{1}+\tau_{0}(s) u_{3}, \quad \frac{d u_{3}}{d s}=-\tau_{0}(s) u_{2} \tag{49}
\end{equation*}
$$

is defined on triples of vectors $\left(u_{1}, u_{2}, u_{3}\right)$. This is the differential equation

$$
\frac{d p_{i}}{d t}=\sum_{j=1}^{3} a_{i j} u_{j}, \quad\left[a_{i j}\right]:=\left[\begin{array}{ccc}
0 & \kappa_{0} & \tau_{0} \\
-\kappa_{0} & 0 & \tau_{0} \\
0 & \tau_{0} & \kappa_{0}
\end{array}\right]
$$

Pick any $c \in I$. Since $\kappa_{0}$ and $\tau_{0}$ are both $C^{r-3}$, and $r \geq 3$, there is a unique $C^{1}$ solution $\left(T_{0}(s), N_{0}(s), B_{0}(s)\right)$ defined on all $I$ which satisfy the initial condition

$$
T_{0}(c)=(1,0,0), \quad N_{0}(c)=(0,1,0), \quad B_{0}(c)=(0,0,1)
$$

Let $\alpha(s):=\int_{c}^{s} T_{0}(s) d s$.
§50. By the existence and uniqueness theorem for differential equations, $T_{0}, N_{0}$, and $B_{0}$ are $C^{r-3}$. From the second equation of (49), $d N_{0} / d s$ is the product of $C^{r-3}$ functions, so $N_{0}$ is $C^{r-2}$. Since $\kappa_{0}$ is $C^{r-2}$ and $N_{0}$ is $C^{r-2}$, the first equation of (49) implies $d T_{0} / d s$ is the product of $C^{r-2}$ functions, so $T_{0}$ is $C^{r-1}$. Since $\alpha^{\prime}=T_{0}$, this shows that $\alpha^{\prime}$ is $C^{r}$.
§51. Define the $3 \times 3$ matrix $\left[p_{i j}\right]:=\left[u_{i} \cdot u_{j}\right]$ and note that

$$
\frac{d p_{i j}}{d t}=\frac{d u_{i}}{d t} \cdot u_{j}+u_{i} \cdot \frac{d u_{j}}{d t}=\left(\sum_{k=1}^{3} a_{i k} u_{k}\right) \cdot u_{j}+u_{i} \cdot\left(\sum_{k=1}^{3} a_{j k} u_{k}\right)
$$

so the $p_{i j}$ satisfy the (linear) initial value problem

$$
\frac{d p_{i j}}{d t}=\sum_{k=1}^{3}\left(a_{i k} p_{k j}+a_{j k} p_{i k}\right), \quad\left[p_{i j}\right]=\mathbf{1}_{3 \times 3}
$$

Constant $p_{i j}$ with the same initial data is also a solution because substitution gives 0 on the left side and $a_{i j}+a_{j i}=0$ on the right, so the vectors $u_{1}, u_{2}, u_{3}$ are orthonormal for all $t$, by uniqueness of initial value problems. The determinant of the $3 \times 3$ matrix $\left[u_{1}, u_{2}, u_{3}\right.$ ] is continuous in $t$ and either +1 or -1 , and is +1 at $t=c$, so it is +1 for all $t$, and $u_{1}, u_{2}, u_{3}$ is a right-handed orthonormal basis.
$\S 52$. Let $(\kappa, \tau, T, N, B)$ be the Serret-Frenet data of $\alpha(s)$, which is unit speed because $\alpha^{\prime}(s)=T_{0}(s)$ and $T_{0}(s)$ is a unit vector for all $s$. Then $T=\alpha^{\prime}=T_{0}$ and $\kappa=\left|T^{\prime}\right|=\left|T_{0}^{\prime}\right|=\left|\kappa_{0} T_{0}\right|=\kappa_{0}$. Also $N=T^{\prime} / \kappa=T_{0}^{\prime} / \kappa_{0}=N_{0}$, and $B=T \times N=T_{0} \times N_{0}=B_{0}$ because $\left(T_{0}, N_{0}, B_{0}\right)$ is right-handed and orthonormal, and $\tau=-B^{\prime} \cdot N=-B_{0}^{\prime} \cdot N_{0}=$ $-\left(-\tau_{0} N_{0}\right) \cdot N_{0}=\tau_{0}$. Thus the Serret-Frent data for $\alpha$ is $\left(\kappa_{0}, \tau_{0}, T_{0}, N_{0}, B_{0}\right)$ and $\alpha$ has the specified curvature and torsion.
$\S 53$. Suppose $\alpha_{1}$ and $\alpha_{2}$ both have curvature $\kappa_{0}$ and torsion $\tau_{0}$. The Serret-Frenet frames $\left(T_{i}, N_{i}, B_{i}\right)$ of $\alpha_{i}, i=1,2$, both satisfy the differential equation (49), possibly with different initial conditions at $s=c$. Let $A$ be the $3 \times 3$ rotation matrix such that

$$
A T_{1}(c)=T_{2}(c), \quad A N_{1}(c)=N_{2}(c), \quad A B_{1}(c)=B_{2}(c)
$$

ie $A=A_{2} A_{1}^{-1}$ where $A_{i}$ are the matrices with column vectors $T_{i}(c), N_{i}(c), B_{i}(c)$. The curve $A \alpha_{1}(s)$ also has curvature $\kappa_{0}$ and torsion $\tau_{0}$, so its Serret-Frenet frame $A \tilde{T}_{1}(s), A \tilde{N}_{1}(s), A \tilde{B}_{1}(s)$ satisfies the same initial value problem as the Serret-Frenet frame of $\alpha_{2}$. By uniqueness of solutions to initial value problems, $A T_{1}(s)=T_{2}(s)$, so

$$
A \alpha_{1}(s)-A \alpha_{1}(c)=\int_{c}^{s} \frac{d}{d s} A \alpha_{1}(s) d s=\int_{c}^{s} A T_{1}(s) d s=\int_{c}^{s} T_{2}(s) d s=\int_{c}^{s} \frac{d}{d s} \alpha_{2}(s) d s=\alpha_{2}(s)-\alpha_{2}(c)
$$

ie $\alpha_{2}(s)=A \alpha_{1}(s)+a$, where $a:=\alpha_{2}(c)-A \alpha_{1}(c)$.
$\S 54$. Remark (Curvature, torsion, and rates of change of angles). Given a unit vector $u(t) \in \mathbb{R}^{3}$, one can ask for the rate of change of the angle $\theta$ that $u(t)$ makes with a given unit vector $\hat{a}$ at $t=a$ : Since $\cos \theta=u(t) \cdot \hat{a}$, differentiate implicitly to obtain

$$
-\sin \theta \frac{d \theta}{d t}=\frac{d u}{d t} \cdot \hat{a}
$$

Choosing $\hat{a}= \pm u^{\prime}(a) /\left|u^{\prime}(a)\right|$ makes sense because then the angle is measured with respect to the same direction as the velocity at $t=p$. For that choice, $\theta(a)=\pi / 2$, because $u(a) \cdot u^{\prime}(a)=(u(a) \cdot u(a))^{\prime}=0$, and $\theta^{\prime}(a)=-u^{\prime}(a) \cdot \hat{a}$. Applying this idea to $T(s)$ and $B(s)$ gives the following:

1. Take $u=T(s), \hat{a}=T^{\prime}(a) /\left|T^{\prime}(a)\right|=\alpha^{\prime \prime}(a) /\left|\alpha^{\prime}(a)\right|=N(a)$, and define $\theta_{\kappa}$ to be the angle between $T(s)$ and $\hat{a}$.

Then $\kappa(a)=\left|\alpha^{\prime \prime}(a)\right|=\left|T^{\prime}(a)\right|$ so $-\kappa(a)=-T^{\prime}(a) \cdot T^{\prime}(a) /\left|T^{\prime}(a)\right|=\theta_{\kappa}^{\prime}(a)$.
2. $B^{\prime}(a)$ is orthogonal to both $T(a)$ and itself since

$$
B^{\prime}(a)=T^{\prime}(a) \times N(a)+T(a) \times N^{\prime}(a)=T^{\prime}(a) \times\left(\kappa(a) T^{\prime}(a)\right)+T(a) \times N^{\prime}(a)=T(a) \times N^{\prime}(a)
$$

So $B^{\prime}(a)$ can only be in the direction of $N(a)$, and one can take $\hat{a}=N(a)$ and define $\theta_{\tau}$ as the angle between $B(s)$ and $N(a)$. Then $\tau(a)=-N(a) \cdot B^{\prime}(a)=-B^{\prime}(s) \cdot N(a)=\theta_{\tau}^{\prime}(a)$.

## §55. Theorem.

1. A line is an immersed curve in $\mathbb{R}^{3}$ which can be reparametrized to $\alpha(t)=a+b t$, where $a$ and $b$ are constant vectors, $b \neq 0$. For $C^{2}$ immersed curves in $\mathbb{R}^{3}$, the follow statements are equivalent: (A) $\alpha$ is a line; ( $B$ ) $T$ is constant; (C) $\kappa=0$.
2. An immersed curve $\alpha(t)$ in $\mathbb{R}^{3}$ is planar if there is a nonzero constant vector $a$ such that $\alpha(t) \cdot a$ is constant. For $C^{3}$ immersed curves $\alpha$ in $\mathbb{R}^{3}$ with $\kappa$ never zero, the following statements are equivalent: (A) $\alpha$ is planar; (B) $B$ is constant; (C) $\tau=0$.
3. An immersed curve $\alpha(t)$ in $\mathbb{R}^{3}$ is a sphere curve (with center a) if there is a constant nonzero vector a such that $|\alpha(t)-a|$ is constant. For $C^{4}$ immersed curves $\alpha$ in $\mathbb{R}^{3}$ with $\kappa$ never zero and $\tau$ never zero, the following statements are equivalent: (A) $\alpha$ is a sphere curve; $(B) \alpha+\rho N+\left(\rho^{\prime} \sigma / v\right) B$ is constant; ( $C$ ) $\left(\sigma \rho^{\prime} / v\right)^{\prime} / v+\rho \tau=0$.
4. A circle is a planar sphere curve. For $C^{3}$ immersed curves $\alpha$ in $\mathbb{R}^{3}$ with $\kappa$ never zero, the following statements are equivalent: (A) $\alpha$ is a circle; (B) $\alpha+\rho N$ is constant and $B$ is constant; ( $C$ ) $\kappa$ is constant and $\tau=0$.
5. A helix (with pitch $|a|$ and axial vector $a$ ) is an immersed curve in $\mathbb{R}^{3}$ which can be reparametrized to $\alpha(t)=\bar{\alpha}(t)+t a$, where $a \neq 0$ is constant and $\bar{\alpha}$ is unit speed and in a plane with normal $a$. For $C^{3}$ immersed curves $\alpha$ in $\mathbb{R}^{3}$ with $\kappa$ never zero and $\tau$ never zero, the following statements are equivalent: (A) $\alpha$ is a helix; (B) there is a nonzero constant vector $a$ such that $T \cdot a$ is a nonzero constant; ( $C$ ) $\tau / \kappa$ is a nonzero constant.
6. A circular helix is an immersed curve in $\mathbb{R}^{3}$ which can be reparametrized to $\alpha(t)=\bar{\alpha}(t)+a t$ where where $a \neq 0$ is constant and $\bar{\alpha}(t)$ is a circle in a plane with normal $a$. For $C^{3}$ immersed curves $\alpha$ in $\mathbb{R}^{3}$ with $\kappa$ never zero and $\tau$ never zero, the following statements are equivalent: (A) $\alpha$ is a circular helix; (B) there is a nonzero vector $a$ such that $T \cdot a$ is a nonzero constant, and $\left(\alpha+\rho v^{2} N\right) \times a$ is constant; ( $C$ ) $\kappa$ and $\tau$ are nonzero constants.
Proof. In the proofs of (1) through (4), assume without loss of generality that $\alpha$ is unit speed.
$(1 \mathrm{~A}) \Rightarrow(1 \mathrm{~B})$ : Given $\alpha=a+b t, \alpha^{\prime}=b$ and $T=\alpha^{\prime} /\left|\alpha^{\prime}\right|=b /|b|$, which is constant.
$(1 \mathrm{~B}) \Rightarrow(1 \mathrm{C})$ : Given that $T$ is constant, $\boldsymbol{v}=v T$, so $\boldsymbol{a}=v^{\prime} T+v T^{\prime}=v^{\prime} T$, and $\boldsymbol{v} \times \boldsymbol{a}=(v T) \times\left(v^{\prime} T\right)=0$ so $\kappa=|\boldsymbol{v} \times \boldsymbol{a}| / v^{3}=0$.
$(1 \mathrm{C}) \Rightarrow(1 \mathrm{~A})$ : Reparametrize $\alpha$ by arc-length. Then $0=\kappa=\left|\alpha^{\prime \prime}\right|=\left|T^{\prime}\right|$ implies $T^{\prime}=0$, so $T$ is constant.
$(2 \mathrm{C}) \Leftrightarrow(2 \mathrm{~B}): \tau=0$ if and only if $B$ is constant, because $B^{\prime}=-\tau N$ and $N \neq 0$.
$(2 \mathrm{~B}) \Rightarrow(2 \mathrm{~A})$ : Given $B$ constant, $(B \cdot \alpha)^{\prime}=B^{\prime} \cdot \alpha+B \cdot \alpha^{\prime}=B \cdot T=0$ so $B \cdot \alpha$ is constant.
$(2 \mathrm{~A}) \Rightarrow(2 \mathrm{~B})$ : Given $\alpha \cdot a=0$, differentiate to obtain $T \cdot a=0$ and $\kappa N \cdot a=0$. So $a$ is parallel to $B$ for all $s$, so $B= \pm a /|a|$ and $B$ is constant because it is continuous.
$(3 \mathrm{~A}) \Rightarrow(3 \mathrm{C})$ : Differentiate the defining property $(\alpha-a) \cdot(\alpha-a)=$ constant as many times as necessary in order to uncover the required functional relation between curvature and torsion:

$$
\begin{aligned}
& (\alpha-a) \cdot(\alpha-a)=\text { constant } \\
& 2 \alpha^{\prime} \cdot(\alpha-a)=0, \quad \text { so } \quad T \cdot(\alpha-a)=0 \\
& T^{\prime} \cdot(\alpha-a)+T \cdot \alpha^{\prime}=\kappa N \cdot(\alpha-a)+1=0, \quad \text { so } \quad N \cdot(\alpha-a)=-\rho, \\
& N^{\prime} \cdot(\alpha-a)+N \cdot \alpha^{\prime}=(-\kappa T+\tau B) \cdot(\alpha-a)=-\rho^{\prime} \quad \text { so } \quad B \cdot(\alpha-a)=-\sigma \rho^{\prime}, \\
& \left(-\sigma \rho^{\prime}\right)^{\prime}=(B \cdot(\alpha-a))^{\prime}=B^{\prime} \cdot(\alpha-a)+B \cdot T=-\tau N \cdot(\alpha-a)=\rho \tau,
\end{aligned}
$$

i.e., $\left(\sigma \rho^{\prime}\right)^{\prime}+\rho \tau=0$.
$(3 \mathrm{C}) \Rightarrow(3 \mathrm{~B})$ : Assuming $\left(\sigma \rho^{\prime}\right)^{\prime}+\rho \tau=0$,

$$
\begin{aligned}
\left(\alpha+\rho N+\left(\rho^{\prime} \sigma\right) B\right)^{\prime} & =\alpha^{\prime}+\rho^{\prime} N+\rho N^{\prime}+\left(\rho^{\prime} \sigma\right)^{\prime} B+\left(\rho^{\prime} \sigma\right) B^{\prime} \\
& =T+\rho^{\prime} N+\rho(-\kappa N+\tau B)+\left(\rho^{\prime} \sigma\right)^{\prime} B+\rho^{\prime} \sigma(-\tau N) \\
& =(1-\kappa \rho) T+\rho^{\prime}(1-\tau \sigma) N+\left(\rho \tau+\left(\sigma \rho^{\prime}+\rho \tau\right)^{\prime}\right) B=0
\end{aligned}
$$

$(3 \mathrm{~B}) \Rightarrow(3 \mathrm{~A})$ : Assuming $\alpha+\rho N+\left(\rho^{\prime} \sigma\right) B=a,((\alpha-a) \cdot(\alpha-a))^{\prime}=2 T \cdot(\alpha-a)=-T \cdot\left(\rho N+\left(\rho^{\prime} \sigma\right) B\right)=0$, so there is a constant distance from $\alpha$ to the constant $a$.
$(4 \mathrm{~A}) \Rightarrow(4 \mathrm{C})$ : Assume $|\alpha-a|$ and $\alpha \cdot b$ are both constant, $b \neq 0$, i.e., assume $\alpha$ is a sphere curve and $\alpha$ is planar, respectively. Without loss of generality, $b$ is a unit vector. Note that any multiple of $b$ can be added to $a$, since

$$
|\alpha-(a+f b)|^{2}=|\alpha|^{2}+2 \alpha \cdot(a+f b)+|a+f b|^{2}=|\alpha+a|^{2}+2 f \alpha \cdot b+|a+f b|^{2}-|a|^{2}
$$

and the left is constant if and only if the right is. Choose $f$ so that $a$ is in the same plane as $\alpha$, i.e., $a \cdot b=\alpha \cdot b$, or equivalently, $(\alpha-a) \cdot b=0$. As in $(3 \mathrm{~A}) \Rightarrow(3 \mathrm{C})$, differentiating $(\alpha-a) \cdot(\alpha-a)=0$ gives $T \cdot(\alpha-a)=0$ and $N \cdot(\alpha-a)=\rho$. Also $(\alpha-a) \cdot b=0$ gives $T \cdot b=N \cdot b=0$ so $B= \pm b, B \cdot(\alpha-a)=0$, and

$$
|\alpha-a|^{2}=(T \cdot(\alpha-a))^{2}+(N \cdot(\alpha-a))^{2}+(B \cdot(\alpha-a))^{2}=\rho^{2}
$$

implying $\kappa=1 / \rho$ is constant.
$(4 \mathrm{C}) \Rightarrow(4 \mathrm{~A})$ : Assume $\alpha(s)$ has constant nonzero curvature $\kappa$ and torsion $\tau=0$. The circle $c(s)$ defined by

$$
x=\frac{1}{\kappa} \cos (\kappa s), \quad y=\frac{1}{\kappa} \sin (\kappa s), \quad z=0, \quad t \in(a, b)
$$

has curvature $\kappa$ and torsion $\tau=0$. By the fundamental theorem for curves, $\alpha=A c(s)+B$ for some $(A, b) \in \mathrm{SE}(3)$, and it follows easily that $\alpha$ is a circle.
$(4 \mathrm{~B}) \Leftrightarrow(4 \mathrm{C}): B^{\prime}=-\tau N$ so $B$ is constant if and only if $\tau=0$, so for either of $(4 \mathrm{~B}) \Rightarrow(4 \mathrm{C})$ or $(4 \mathrm{~B}) \Leftarrow(4 \mathrm{C})$ one can assume $\tau=0$ and $B$ is constant. Then $(\alpha+\rho N)^{\prime}=T+\rho N^{\prime}+\rho^{\prime} N=T+\rho(-\kappa T+\tau B)+\rho^{\prime} N=\rho^{\prime} N$, which shows $\alpha+\rho N$ is constant if and only if $\rho$ is constant, and $\kappa=1 / \rho$.
$(5 \mathrm{~A}) \Rightarrow(5 \mathrm{~B})$ : Differentiating $\alpha(t)=\bar{\alpha}(t)+$ at gives $\alpha^{\prime}=\bar{\alpha}^{\prime}+a$ where $\left|\bar{\alpha}^{\prime}\right|=1$ and $\bar{\alpha}^{\prime} \cdot a=0$. Thus $\left|\alpha^{\prime}\right|^{2}=1+|a|^{2}$ and

$$
T \cdot a=\frac{1}{\left|\alpha^{\prime}\right|} \alpha^{\prime} \cdot a=\frac{1}{\sqrt{1+|a|^{2}}}\left(\bar{\alpha}^{\prime}+a\right) \cdot a=\frac{|a|^{2}}{\sqrt{1+|a|^{2}}}
$$

which is constant.
$(5 \mathrm{~B}) \Rightarrow(5 \mathrm{~A})$ : Without loss of generality, $\alpha$ is unit speed and $a$ is a unit vector. Let $h$ be the constant $T \cdot a$ and define $\bar{\alpha}(s)=\alpha(s)-(h a) s$. Then $\alpha^{\prime}=\bar{\alpha}^{\prime}+h a$, so $\bar{\alpha}^{\prime} \cdot a=\alpha^{\prime} \cdot a-h a \cdot a=h-h=0$, and $1=\left|\alpha^{\prime}\right|^{2}=\left|\bar{\alpha}^{\prime}\right|^{2}+h^{2}$, i.e., $\left|\bar{\alpha}^{\prime}\right|=$ $\sqrt{1-h^{2}}$, a constant. If $\bar{\alpha}^{\prime}(s)$ is ever zero, then it is zero for all $s$, and $\alpha(s)=(h a) s$, contradicting the assumption that the curvature of $\alpha$ is never zero. Thus $\bar{\alpha}$ is immersed and has constant, nonzero speed. Integrating $\bar{\alpha}^{\prime}(s) \cdot a=0$ implies $\bar{\alpha} \cdot a$ is constant, so $\bar{\alpha}(s)$ is planar. Reparametrize $\alpha(s)$ and $\bar{\alpha}(s)$ by $t=\sqrt{1-h^{2}} s$, so $\alpha(t) \equiv \alpha\left(t / \sqrt{1-h^{2}}\right)$ and $\bar{\alpha}(t) \equiv \bar{\alpha}\left(t / \sqrt{1-h^{2}}\right)$. Substituting into $\alpha(s)=\bar{\alpha}(s)+s(h a)$ gives $\alpha(t)=\bar{\alpha}(t)+\left(h / \sqrt{1-h^{2}} a\right) t$ and $\bar{\alpha}(t)$ is unit speed in a plane with normal $a$.
$(5 \mathrm{~B}) \Rightarrow(5 \mathrm{C})$ : Without loss of generality, assume $\alpha(s)$ is unit speed. Differentiating $T \cdot a$ constant gives $T^{\prime} \cdot a=0$, so $N \cdot a=\left(T^{\prime} / \kappa\right) \cdot a=0, N^{\prime} \cdot a=0$, and $0=N^{\prime} \cdot a=(-\kappa T+\tau B) \cdot a=-(T \cdot a) \kappa+(B \cdot a) \tau$. This implies $\kappa$ and $\tau$ have constant ratio because $T \cdot a$ is constant, and $B \cdot a$ is constant follows from $(B \cdot a)^{\prime}=B^{\prime} \cdot a=-\tau N \cdot a=0$.
$(5 \mathrm{C}) \Rightarrow(5 \mathrm{~B})$ : Let $h=\tau / \kappa$, a constant, and define $a=h T+B$. Then $a^{\prime}=h T^{\prime}+B^{\prime}=h(\kappa N)+(-\tau N)=(h \kappa-\tau) N=0$ so $a$ is constant and $T \cdot a=T \cdot(h T+B)=h$.
$(6 \mathrm{~A}) \Rightarrow(6 \mathrm{~B})$ : A circular helix is a helix so, by $(5), T \cdot a$ is constant. Assuming $\alpha(t)=\bar{\alpha}(t)+a t$ where $a \cdot \alpha$ is constant and $\bar{\alpha}$ is a circle in a plane orthogonal to $a$, it follows that $\alpha^{\prime}=\bar{\alpha}^{\prime}+a$ and $\bar{\alpha}^{\prime} \cdot a=0$ so $\left|\alpha^{\prime}\right|^{2}=1+|a|^{2}$ and $v=\left|\alpha^{\prime}\right|$ is constant. So $\boldsymbol{a}=\alpha^{\prime \prime}=\kappa v^{2} N=\bar{\alpha}^{\prime \prime}=\bar{\kappa} \bar{N}$, so $\kappa v^{2}=\bar{\kappa}$ and $N=\bar{N}$, where the bars indicate the Serret-Frenet data of $\bar{\alpha}$. Thus $a \times\left(\alpha+\rho v^{2} N\right)=a \times\left(\bar{\alpha}+a t+\rho / v^{2} N\right)=a \times(\bar{\alpha}+\bar{\rho} \bar{N}+a t)=a \times(\bar{\alpha}+\bar{\rho} \bar{N})$ which is constant since $\bar{\alpha}$ is a circle.
$(6 \mathrm{~B}) \Rightarrow(6 \mathrm{~A})$ : If $T \cdot a$ is constant then by (5) $\alpha$ is a helix, and $\alpha(t)=\bar{\alpha}(t)+a t$ where $\bar{\alpha}$ is in a plane with normal $a$. As in $(6 \mathrm{~A}) \Rightarrow(6 \mathrm{~B}), v=\left|\alpha^{\prime}\right|$ is constant, $\kappa v^{2}=\bar{\kappa}$, and $N=\bar{N}$. Since $a \times\left(\alpha+\rho / v^{2} N\right)$ is constant, it follows that $\alpha+\rho / v^{2} N=c+f(t) a$ where $c$ is constant and $f(t)$ is a function of $t$. Then $\bar{\alpha}=\alpha-a t=c+f a-\rho / v^{2} N-a t=$ $c+f a-\bar{\rho} \bar{N}-a t$, and $a \cdot \bar{\alpha}=a \cdot c+(f-t)|a|^{2}$ is constant since $\alpha$ in in a plane with normal $a$. Thus $f=t$ and $\bar{\alpha}=c-\bar{\rho} \bar{N}$, i.e., $\bar{\alpha}+\bar{\rho} \bar{N}$ is constant.
$(6 \mathrm{~A}) \Rightarrow(6 \mathrm{C})$ : As in $(6 \mathrm{~A}) \Rightarrow(6 \mathrm{~B})$, from $\alpha=\bar{\alpha}+$ at follows $v=\left|\alpha^{\prime}\right|$ is constant and $\kappa v^{2}=\bar{\kappa}$. Thus $\kappa$ is constant, and $\tau / \kappa$ is constant by (5), so $\tau$ is also constant.
$(6 \mathrm{C}) \Rightarrow(6 \mathrm{~A})$ : If $\kappa$ and $\tau$ are constant then so is $\tau / \kappa$, so $\alpha$ is a helix which can be parametrized as $\alpha(t)=\bar{\alpha}(t)+a t$, as in $(5)$. As in $(6 \mathrm{~A}) \Rightarrow(6 \mathrm{~B}), v$ is constant and $\kappa v^{2}=\bar{\kappa}$, so $\bar{\kappa}$ is constant. Thus $\bar{\alpha}$ is a circle.
$\S 56$. Theorem. For $C^{3}$ immersed curve $\alpha(t)$ (not necessarily unit speed) such that $\alpha^{\prime \prime}(t) \neq 0$ for all $t$,

$$
\boldsymbol{v}=v T, \quad \boldsymbol{a}=\frac{d v}{d t} T+v^{2} \kappa N, \quad T=\frac{\boldsymbol{v}}{v}, \quad B=\frac{\boldsymbol{v} \times \boldsymbol{a}}{|\boldsymbol{v} \times \boldsymbol{a}|}, \quad N=B \times T, \quad \kappa=\frac{|\boldsymbol{v} \times \boldsymbol{a}|}{v^{3}}, \quad \tau=\frac{(\boldsymbol{v} \times \boldsymbol{a}) \cdot \boldsymbol{a}^{\prime}}{|\boldsymbol{v} \times \boldsymbol{a}|^{2}},
$$

where $(\kappa, \tau, T, N, B)$ is the Serret-Frenet data of any arc-length parameterization.
Proof.

$$
\begin{aligned}
& \frac{d s}{d t}=\frac{d}{d t} \int_{p}^{t}\left|\frac{d \alpha}{d t}\right| d t=\left|\frac{d \alpha}{d t}\right|=v(t) \\
& \boldsymbol{v}=\frac{d \alpha}{d t}=\frac{d \alpha}{d s} \frac{d s}{d t}=v T
\end{aligned}
$$

$$
\begin{aligned}
& \boldsymbol{a}=\frac{d \boldsymbol{v}}{d t}=\frac{d}{d t}(v T)=\frac{d v}{d t} T+v \frac{d T}{d s} \frac{d s}{d t}=\frac{d v}{d t} T+\kappa v^{2} N \\
& \boldsymbol{v} \times \boldsymbol{a}=(v T) \times\left(\frac{d v}{d t} T+\kappa v^{2} N\right)=\kappa v^{3} T \times N=\kappa v^{3} B \\
& |\boldsymbol{v} \times \boldsymbol{a}|=\left|\kappa v^{3} B\right|=\kappa v^{3}, \quad \text { so } \quad \kappa=\frac{|\boldsymbol{v} \times \boldsymbol{a}|}{v^{3}} \quad \text { and } \quad B=\frac{\boldsymbol{v} \times \boldsymbol{a}}{\kappa v^{3}}=\frac{\boldsymbol{v} \times \boldsymbol{a}}{|\boldsymbol{v} \times \boldsymbol{a}|}, \\
& B \times T=-T \times(T \times N)=-(T \cdot N) T+(T \cdot T) N=N, \\
& (\boldsymbol{v} \times \boldsymbol{a}) \cdot \frac{d \boldsymbol{a}}{d t}=|v \times a| B \cdot \frac{d}{d t}\left(\frac{d v}{d t} T+\kappa v^{2} N\right)=|v \times a| B \cdot\left(\frac{d^{2} v}{d t^{2}} T+\frac{d v}{d t} \kappa N+\frac{d}{d t}\left(\kappa v^{2}\right) N+\kappa v^{2} \frac{d N}{d t}\right)=|v \times a| \kappa v^{3} \tau, \\
& \tau=\frac{\boldsymbol{v} \times \boldsymbol{a}}{|v \times a| \kappa v^{3}}=\frac{\boldsymbol{v} \times \boldsymbol{a}}{|v \times a|^{2}}
\end{aligned}
$$

## 57. Foundations

§58. Differential geometry is a generalization of multivariate calculus. Some essential elements of multivariate calculus are summarized here.
§59. Many central results of differential geometry are applications of at least one of three fundamental theorems:

1. the inverse function theorem, which provides smooth solutions to general equations, and a plentiful source of smoothly invertible maps;
2. the existence and uniqueness theorem for ordinary differential equations, which provides smooth unique solutions for the initial value problems of systems of ordinary differential equations; and
3. the Frobenius theorem, which, generalizing the computation of a potential from a conservative vector field, provides solutions to a particular kind of system of partial differential equations.
It is necessary to acquire a working knowledge of the underlying multivariate calculus. The details of the proof of the three theorems need not be mastered - these have compelling informal explanations.

## 60. Euclidean space, and the vector space $\mathbb{R}^{n}$

$\S 61$. Euclidean space $\mathbb{R}^{n}$ : By $\mathbb{R}^{n}$ is meant the set of $n$-tuples $x=\left(x_{1}, x_{2}, \ldots, x_{n}\right)$, where $x_{i} \in \mathbb{R}$. The tuples are thought of as places. No particular meaning can generally be assigned to the sum of two places eg to the sum of map coordinates of, say, Los Angeles and New York, so addition of $n$-tuples is undefined when they are thought this way. In this context, $\mathbb{R}^{n}$ is a metric space: the distance between two $n$-tuples $x$ and $y$ is $\sqrt{\left(x_{1}-y_{1}\right)^{2}+\cdots+\left(x_{n}-y_{n}\right)^{2}}$. $\S 62$. The vector space $\mathbb{R}^{n}: \mathbb{R}^{n}$ also denotes the vector space of $n$-tuples with componentwise vector addition and scalar multiplication. Here the $n$-tuples are written as column vectors, so that the usual matrix multiplication can be applied to them with the matrix on the left. In this context $\mathbb{R}^{n}$ is an inner product space, with inner product $v \cdot w:=v^{T} w=v_{1} w_{1}+\cdots v_{n} w_{n}$. The length of a vector $v$ is $|v|:=\sqrt{v \cdot v}$. The Cauchy-Schwartz inequality states $|v \cdot w| \leq|v||w|$; the triangle inequality states that $|v+w| \leq|v|+|w|$. We will make use of the following fact: if $v_{1}, \ldots, v_{n} \in \mathbb{R}^{n}$ are pairwise orthonormal unit vectors, i.e., $v_{i} \cdot v_{j}$ is 1 if $i=j$ and 0 otherwise, then any vector $x \in \mathbb{R}^{n}$ is uniquely a sum of the $v_{i}$ by the formula

$$
\begin{equation*}
x=\left(v_{1} \cdot x\right) v_{1}+\left(v_{2} \cdot x\right) v_{2}+\cdots+\left(v_{n} \cdot x\right) v_{n} \tag{63}
\end{equation*}
$$

Such a tuple $\left(v_{1}, \ldots, v_{n}\right)$ is called an orthonormal basis.
$\S 64$. The Cauchy-Schwartz inequality follows by expanding $(t v+w) \cdot(t v+w) \geq 0$ and putting $t=-v \cdot w /|v|^{2}$ (which is the value of $t$ for which that expression has a global minimum). The triangle inequality follows from Cauchy-Schwartz via $|v+w|^{2}=(v+w) \cdot(v+w)=|v|^{2}+2 v \cdot w+|w|^{2}$. Pairwise orthonormal vectors $v_{i}$ are linearly independent: dotting $c_{1} v_{1}+c_{2} v_{2}+\cdots+c_{n} v_{n}=0$ with $v_{i}$ gives $c_{1} v_{i} \cdot v_{1}+c_{2} v_{i} \cdot v_{2}+\cdots+c_{n} v_{i} \cdot v_{n}=c_{i}=0$. $n$ linearly independent vectors in $\mathbb{R}^{n}$ form a basis: any vector is uniquely a linear combination of them. Dotting $x=a_{1} v_{1}+a_{2} v_{2}+\cdots+a_{n} v_{n}$ with $v_{i}$ gives $v_{i} \cdot x=a_{1} v_{i} \cdot v_{1}+a_{2} v_{i} \cdot v_{2}+\cdots+a_{n} v_{i} \cdot v_{n}=a_{i}$.
$\S 65$. The basic constructs in differential geometry - curves, surfaces, manifolds, and submanifolds - are locally Euclidean. Linearization is an essential theme, the natural carrier of which is a vector space. Every vector space is the same as $\mathbb{R}^{n}$ after choosing a basis. The triangle inequality is an often used tool to bound a sum of vectors in terms of the summands.
§66. Summing places is nonsensical, but the difference of two places may not be: it does make sense to regard the difference of places $x$ and $y$ in $\mathbb{R}^{n}$ as the column vector $v$, where $v_{i}=y_{i}-x_{i}$. That is the vector from $x$ to $y$ and it is typically visualized as an arrow with tail at $x$ and head at $y$. The distance between $x$ and $y$ is the length of the vector $v$. It also makes sense to sum a vector $v$ and a place $x$ to obtain the new place $y=x+v$. The vectors like $v$ are steps between the places. While the difference between places and vectors is conceptually significant, the sets of vectors and places have the same name $\mathbb{R}^{n}$. Vectors ought always be written as a column but are often written as a row in order to avoid awkwardly using up vertical space on the page. A linear mapping $\alpha$ on vectors with values in $\mathbb{R}$ is regarded as a row vector, so that its value on a vector $v$ is the matrix product $\alpha v$. The set of such dual vectors is a vector space which is conceptually different from the metric space of places, but those are both $n$-tuples written as rows. Differential geometry is conceptual. You have to get used to carrying the burden of intepretation. A tuple of real numbers may represent conceptually different objects, and you have to fully understand which tuple is representing which kind of object.
$\S 67$. Open and closed sets; continuous functions: The open balls and closed balls around $x \in \mathbb{R}^{n}$ of radius $r$ are respectively

$$
B_{r}(x):=\left\{y \in \mathbb{R}^{n}| | y-x \mid<r\right\}, \quad \bar{B}_{r}(x):=\left\{y \in \mathbb{R}^{n}| | y-x \mid \leq r\right\}
$$

$A \subseteq \mathbb{R}^{n}$ is an open subset of $\mathbb{R}^{n}$ if it contains an open ball around each of its points i.e. if for all $x \in A$ there is an $r$ such that $B_{r}(x) \subseteq A . K \subseteq \mathbb{R}^{n}$ is a closed subset if $\mathbb{R}^{n} \backslash K$ is open. A function $f: A \rightarrow \mathbb{R}^{m}$ is continuous at $a$ if $|f(x)-f(a)|$ is arbitrarily small whenever $|x-a|$ is sufficiently small; $f$ is continuous if it continuous at all $a$. Some useful facts:

1. If $f: A \rightarrow \mathbb{R}^{n}$ is continuous, and $A$ is open, then $f^{-1}(V):=\{x \mid f(x) \in V\}$ is open for every open set $V \subseteq \mathbb{R}^{n}$. If $f: A \rightarrow \mathbb{R}^{n}$ is continuous and $A$ is closed, then $f^{-1}(K):=\{x \mid f(x) \in K\}$ is closed for every closed set $K \subseteq \mathbb{R}^{n}$.
2. $K$ is a closed subset of $\mathbb{R}^{n}$ if and only if all convergent sequences in $K$ converge to an element of $K$, i.e., if $x_{i} \in K$ and $\lim _{i \rightarrow \infty} x_{i}=x$ then $x \in K$.
§68. Open sets encapsulate membership stability: if $A$ is open, $a \in A$, and $a$ is moved a little, then still $a \in A$. Also, it is often important that a point in a set is surrounded by other points in the same set. For example, the derivative of a function $f(x)$ at $x=a$ is defined as

$$
\frac{d f}{d x}=\lim _{h \rightarrow 0} \frac{f(a+h)-f(a)}{h}
$$

where the limit means that for all $\epsilon>0$ there is a $\delta>0$ such that $|f(a+h)-f(a)|<\epsilon$ whenever $x \in \operatorname{domain} f$ and $|h|<\delta$. If there is a disk about $x=a$ which meets the domain of $f$ only at $x=a$ then the limit definition is vacuously satisfied for any value of the limit and the derivative is not unique. So, the usual expositions restrict the definition of the partial derivative to functions defined only on open intervals where the limit is unique if it exists at all, rather than have to postulate such uniqueness in the definition of the derivative itself. Closed sets are often used when it is required to arrange that the limit of a sequence in a set is also in that set.
§69. These basic topological notions and results may be found in standard advanced calculus texts, such as [5], [6], [7]. The proofs are easy: For the first fact, if $a \in f^{-1}(V)$ then $f(a) \in V$ and $V$ is open, so there is a $\epsilon>0$ such that $B_{\epsilon}(f(a)) \subset V$. By continuity of $f$, there is a $\delta_{1}>0$ such that $|f(x)-f(a)|<\epsilon$ whenever $|x-a|<\delta_{1}$, and a $\delta_{2}$ such that $B_{\delta_{2}}(a) \subset A$. Setting $\delta=\min \left(\delta_{1}, \delta_{2}\right)$, one has $B_{\delta}(a) \subseteq A \cap f^{-1}(V)$. For the second fact, if $K$ is closed, and $x_{i} \in K$ and $\lim _{i \rightarrow \infty} x_{i}=x$, then $x \neq K$ is a contradiction because then $\mathbb{R}^{n} \backslash K$ is open and $x_{i}$ could not be within some radius of $x$.
§70. Isometries: An affine map of $\mathbb{R}^{3}$ is a linear map followed by a translation, i.e., a map of the form $x \rightarrow A x+b$ where $A$ is a $3 \times 3$ matrix and $b$ is a vector. An affine map $A x+b$ is a Euclidean isometry if it preserves distance, i.e., if $|(A x+b)-(A y+b)|=|A(x-y)|=|x-y|$. We may have use of the following:

1. A tuple $\left(v_{1}, v_{2}, v_{3}\right)$ of vectors in $\mathbb{R}^{3}$ is right-handed if the $3 \times 3$ matrix $\left[v_{1} v_{2} v_{3}\right.$ ] with columns $v_{1}, v_{2}, v_{3}$ has positive determinant, and left-handed if it is negative. There are exactly three kinds of such tuples: lefthanded, right-handed, and, where the determinant is zero, degenerate. An affine map $A x+b$ is proper if it
preserves handedness, i.e., $\left(A v_{1}, A v_{2}, A v_{3}\right)$ is right handed whenever $\left(v_{1}, v_{2}, v_{3}\right)$ is. The set of proper Euclidean isometries of $\mathbb{R}^{3}$ is called the special euclidean group and it is denoted $\mathrm{SE}(3)$.
2. An affine map $A x+b$ is proper if and only if $\operatorname{det} A>0$. It is a Euclidean isometry if and only if $A^{T} A=\mathbf{1}$, where $\mathbf{1}$ is the identity matrix. It is a proper Euclidean isometry if and only if $A^{T} A=\mathbf{1}$ and $\operatorname{det} A=1$.
3. Given a matrix $A$ such that $\operatorname{det} A=1$ and $A^{T} A=\mathbf{1}, A$ is a rotation about some axis, because there is a right handed orthonormal basis of $\mathbb{R}^{3}$ with respect to which $A$ is

$$
\left[\begin{array}{ccc}
\cos \theta & -\sin \theta & 0 \\
\sin \theta & \cos \theta & 0 \\
0 & 0 & 1
\end{array}\right] \quad \text { or } \quad\left[\begin{array}{ccc}
\cos \theta & \sin \theta & 0 \\
\sin \theta & -\cos \theta & 0 \\
0 & 0 & -1
\end{array}\right],
$$

i.e., either rotations about an axis that preserve their axis or reverse it, respectively.
$\S 71$. All of these results are elementary linear algebra, but they are specialized enough that the proofs are collected into Proposition 72. Note that $A^{T} A=\mathbf{1}$ implies $A^{T}=A^{-1}$ and hence $A A^{T}=A A^{-1}=\mathbf{1}$, and also note that $A^{T} A=1$ is equivalent to $(A u) \cdot(A v)=u \cdot v$ of all $u$ and $v$.

## §72. Proposition.

1. An affine map $A x+b$ is proper if and only if $\operatorname{det} A>0$. It is a Euclidean isometry if and only if $A^{T} A=\mathbf{1}$, where $\mathbf{1}$ is the identity matrix. It is a proper Euclidean isometry if and only if $A^{T} A=\mathbf{1}$ and $\operatorname{det} A=1$.
2. For a $3 \times 3$ matrix $A$, $\operatorname{det} A=1$ and $A^{T} A=\mathbf{1}$ if and only if there is an orthonormal basis with respect to which $A$ is one of

$$
\left[\begin{array}{ccc}
\cos \theta & -\sin \theta & 0 \\
\sin \theta & \cos \theta & 0 \\
0 & 0 & 1
\end{array}\right] \quad \text { or } \quad\left[\begin{array}{ccc}
\cos \theta & \sin \theta & 0 \\
\sin \theta & -\cos \theta & 0 \\
0 & 0 & -1
\end{array}\right]
$$

## Proof.

$\S 73$. For $(1)$, if $\left(v_{1}, v_{2}, v_{3}\right)$ is a right-handed tuple then

$$
\operatorname{det}\left(\left[A v_{1}, A v_{2}, A v_{3}\right]\right)=\operatorname{det}\left(A\left[v_{1}, v_{2}, v_{3}\right]\right)=\operatorname{det}(A) \operatorname{det}\left(\left[v_{1}, v_{2}, v_{3}\right]\right)
$$

so $A x+b$ is a proper affine map if and only if $\operatorname{det} A>0$. If $A x+b$ is an isometry then the distance between 0 and any $u \in \mathbb{R}^{3}$ is preserved, so

$$
|u|^{2}=|u-0|^{2}=|(A u+b)-(A 0+b)|^{2}=|A u|^{2}
$$

and $u^{T} u=|A u|^{2}=(A u)^{T}(A u)=u^{T}\left(A^{T} A\right) u$ for all $u$. It follows that, if $u, v \in \mathbb{R}^{3}$ then

$$
\begin{aligned}
0 & =(u+v)^{T}\left(A^{T} A\right)(u+v)-(u+v)^{T}(u+v) \\
& =\left(u^{T}\left(A^{T} A\right) u+u^{T}\left(A^{T} A\right) v+v^{T}\left(A^{T} A\right) u+v^{T}\left(A^{T} A\right) v\right)-\left(u^{T} u+u^{T} v+v^{T} u+v^{T} v\right) \\
& =2 u^{T}\left(A^{T} A-\mathbf{1}\right) v
\end{aligned}
$$

Setting $u$ and $v$ to the coordinate vectors $\boldsymbol{e}_{i}$ and $\boldsymbol{e}_{j}$, respectively, obtains that the general $(i, j)$ element of $A^{T} A-\mathbf{1}$ is zero, so $A^{T} A=\mathbf{1}$. The converse, that if $A^{T} A=\mathbf{1}$ then $A x+b$ is a Euclidean isometry, is a direct calculation. Also, if $A^{T} A=\mathbf{1}$, then $\operatorname{det}(A)^{2}=\operatorname{det}\left(A^{T} A\right)=\operatorname{det} \mathbf{1}=1$, so $\operatorname{det} A>0$ if and only if $\operatorname{det} A=1$.
$\S 74$. For (2), the characteristic polynomial $\operatorname{det}(A-\lambda \mathbf{1})$ is cubic and so has at least one real root $\lambda . \lambda=0$ is not possible because $\operatorname{det} A=1$. Thus there is a unit vector, say $u_{3}$, such that $A u_{3}=\lambda u_{3}$, and $\left|u_{3}\right|^{2}=\left|A u_{3}\right|^{2}=\lambda^{2}\left|u_{3}\right|^{2}$ so $\lambda= \pm 1$. Let $u_{1}$ and $u_{2}$ be such that $\left(u_{1}, u_{2}, u_{3}\right)$ is a right handed orthonormal basis. From (63),

$$
\begin{aligned}
& A u_{1}=\left(u_{1} \cdot\left(A u_{1}\right)\right) u_{1}+\left(u_{2} \cdot\left(A u_{1}\right)\right) u_{2}+\left(u_{3} \cdot\left(A u_{1}\right)\right) u_{3} \\
& A u_{2}=\left(u_{1} \cdot\left(A u_{2}\right)\right) u_{1}+\left(u_{2} \cdot\left(A u_{2}\right)\right) u_{2}+\left(u_{3} \cdot\left(A u_{2}\right)\right) u_{3} \\
& A u_{3}=\left(u_{1} \cdot\left(A u_{3}\right)\right) u_{1}+\left(u_{2} \cdot\left(A u_{3}\right)\right) u_{2}+\left(u_{3} \cdot\left(A u_{3}\right)\right) u_{3}
\end{aligned}
$$

so the matrix of $A$ with respect to $\left(u_{1}, u_{2}, u_{3}\right)$ is

$$
B=\left[\begin{array}{lll}
u_{1} \cdot\left(A u_{1}\right) & u_{1} \cdot\left(A u_{2}\right) & u_{1} \cdot\left(A u_{3}\right) \\
u_{2} \cdot\left(A u_{1}\right) & u_{2} \cdot\left(A u_{2}\right) & u_{2} \cdot\left(A u_{3}\right) \\
u_{3} \cdot\left(A u_{1}\right) & u_{3} \cdot\left(A u_{2}\right) & u_{3} \cdot\left(A u_{3}\right)
\end{array}\right]^{T}=\left[u_{1}, u_{2}, u_{3}\right] A\left[u_{1}, u_{2}, u_{3}\right]
$$

Now,

$$
u_{3} \cdot\left(A u_{1}\right)= \pm\left(A u_{3}\right) \cdot\left(A u_{1}\right)=u_{3} \cdot u_{1}=0, \quad u_{1} \cdot\left(A u_{3}\right)= \pm u_{1} \cdot\left( \pm u_{3}\right)=0
$$

and similarly $u_{3} \cdot\left(A u_{2}\right)=0$ and $u_{2} \cdot\left(A u_{3}\right)=0$, so the matrix of $A$ with respect to $\left(u_{1}, u_{2}, u_{3}\right)$ has the form

$$
B=\left[\begin{array}{ccc}
a & b & 0 \\
c & d & 0 \\
0 & 0 & \pm 1
\end{array}\right]
$$

The matrix $U:=\left[u_{1}, u_{2}, u_{3}\right]$ also satisfies $U^{T} U=1$ and $\operatorname{det} U=1$, and so $B=U A U^{T}$ does that as well. It follows that $a^{2}+b^{2}=c^{2}+d^{2}=1$ and $a b+c d=0$. Choose $\theta$ so that $a=\cos \theta$ and $c=\sin \theta$. Since the (unit) 2-vector $(b, d)$ is orthogonal to $(a, c)$, there are the two possibilities

$$
(b, d)=(-\sin \theta, \cos \theta) \quad \text { or } \quad(b, d)=(\sin \theta,-\cos \theta)
$$

giving the four possibilities

$$
\left[\begin{array}{ccc}
\cos \theta & -\sin \theta & 0 \\
\sin \theta & \cos \theta & 0 \\
0 & 0 & 1
\end{array}\right], \quad\left[\begin{array}{ccc}
\cos \theta & \sin \theta & 0 \\
\sin \theta & -\cos \theta & 0 \\
0 & 0 & -1
\end{array}\right], \quad\left[\begin{array}{ccc}
\cos \theta & -\sin \theta & 0 \\
\sin \theta & \cos \theta & 0 \\
0 & 0 & -1
\end{array}\right], \quad\left[\begin{array}{ccc}
\cos \theta & \sin \theta & 0 \\
\sin \theta & -\cos \theta & 0 \\
0 & 0 & 1
\end{array}\right]
$$

The third and fourth are not possible because their determinants are -1 .
$\S 75$. Euclidean isometries are everywhere: every time you go somewhere, your new body is position as approximately a proper Euclidean isometry applied to your prior one. They are also the foundation of Euclidean geometry: two triangles are congruent if and only if they can be mapped to one another using a proper Euclidean isometry.

## 76. Multivariate calculus

§77. $C^{r}$ functions: The function $y: A \subseteq \mathbb{R}^{m} \rightarrow \mathbb{R}^{n}$ given by the vector of real valued component functions

$$
y_{1}=y_{1}\left(x_{1}, \ldots, x_{m}\right), \ldots, y_{n}=y_{n}\left(x_{1}, \ldots, x_{m}\right)
$$

is called $C^{r}, r \geq 1$, if $A$ is open and the partial derivatives of the component functions exist and are continuous up to and including order $r$. $C^{0}$ is synonymous with continuous, and $C^{\infty}$ means $C^{r}$ for all $r \geq 0$. A $C^{r}$ diffeomorphism is a $C^{r}$ map with a $C^{r}$ inverse. Dropping the assumption that $A$ is open, a homeomorphism is any continuous map with a continuous inverse.
§78. Derivatives: The derivative of $y: A \subseteq \mathbb{R}^{m} \rightarrow \mathbb{R}^{n}$ is the $n \times m$ matrix $D y \equiv\left[\frac{\partial y_{i}}{\partial x_{j}}\right]$. If $z(y)$ is is another function

$$
z_{1}=z_{1}\left(y_{1}, \ldots, y_{n}\right), \ldots, z_{p}=z_{p}\left(y_{1}, \ldots, y_{n}\right)
$$

and both $z$ and $y$ are $C^{r}$, then, by the chain rule, the result of substituting $y$ into $z$ is also $C^{r}$, and

$$
\frac{\partial z_{k}}{\partial x_{i}}=\frac{\partial z_{k}}{\partial y_{1}} \frac{\partial y_{1}}{\partial x_{i}}+\cdots+\frac{\partial z_{k}}{\partial y_{n}} \frac{\partial y_{n}}{\partial x_{i}}
$$

where the meaning is that equality is obtained after substituting the $y$ in the functions $\partial z_{k} / \partial y_{j}$.
$\S 79$. Any function from $A \subseteq \mathbb{R}^{m}$ to $\mathbb{R}^{m}$ defined by a rule consisting of the operations,,$+- \times, \div$, and substitution, and the including functions $a^{x}, \log _{a} x$, trigonometric functions, power series, and any other $C^{r}$ function, is continuous on its domain and $C^{r}$ on any open subset of $\mathbb{R}^{m}$ contained in its domain. Care has to be exercised with functions
such as $\sqrt{x^{2}}$, the rule for which consists of perfectly ordinary operations, but which is not differentiable at $x=0$ (it is equal to $|x|$, the classic example of a non-differentiable continuous function). To understand this, remember that $x^{a}$ is defined as $e^{a \ln x}$, so $\sqrt{x^{2}}=e^{\frac{1}{2} \ln (x \times x)}$ which is $C^{\infty}$ on any open subset on which it is defined, and so is $C^{\infty}$ on $\{x \mid x \neq 0\}$. The differentiability of $\sqrt{x^{2}}$ at $x=0$ is not assured by the formula $\sqrt{x^{2}}=e^{\frac{1}{2} \ln (x \times x)}$ since that is not defined at $x=0$. Of course, with any of its common meanings for non-integral $a, x^{a}$ is continuous on its domain of definition, $C^{\infty}$ except at $x=0$, and, for $a>0, C^{k}$ at $x=0$, where $k=\lfloor a\rfloor$, if it is defined for both positive and negative $x$, where $k=\lfloor a\rfloor$.
$\S 80$. If $f(x, y) \equiv x y / \sqrt{x^{2}+y^{2}}$ for $(x, y) \neq 0$, and $f(0,0) \equiv 0$, then $\partial f / \partial x$ and $\partial f / \partial y$ exist at all $(x, y)$ but are not continuous at $(0,0)$; and also, if the graph is rotated by $45^{\circ}$ then the partial derivatives of the rotated graph do not exist at $(0,0)$ (along the line $y=x$ the function is $f(x, x)=x^{2} / \sqrt{2 x^{2}}=|x| / \sqrt{2}$ ). So mere existence of partial derivatives is coordinate dependent. On the other hand, composition of a $C^{1}$ function with a $C^{1}$ diffeomorphism is also $C^{1}$ (so in particular composition with a linear change of coordinates, such as a rotation, does not destroy the $C^{1}$ property). The $C^{r}$ property is independent of coordinates: if a function has that property in one coordinate system, then it has that same property in all coordinate systems.
$\S 81$. The class of $C^{r}$ functions provides a convenient framework for many results; for example, a standard result in multivariate calculus is that $C^{1}$ functions are differentiable, in the sense that they are well approximated by their first order Taylor polynomials. Also, if a function is $C^{r}$ and $0 \leq k \leq r$, then any partial derivative of order $k$ is independent of the order of differentiation. The chain rule is needed to compute derivatives through substitutions, something which happens alot.
§82. Taylor expansions, big and little oh notation: Let $f(x)$ be a $C^{r}$ function, $r \geq 0$. The order $r$ Taylor formula is

$$
f(x+h)=f(x)+\frac{f^{(1)}(x)}{1!} h+\cdots+\frac{f^{(r)}(x)}{r!} h^{r}+R_{r}(x, h) h^{r}, \quad R_{r}(x, h) \equiv \int_{0}^{1} \frac{(1-t)^{r-1}}{(r-1)!}\left(f^{(r)}(x+t h)-f^{(r)}(x)\right) d t .
$$

Often, the detailed functional form of the remainder is not of much interest. Since $R_{r}(x, h) h^{r} / h^{r} \rightarrow 0$ as $h \rightarrow 0$, Taylor's formula may be written

$$
f(x+h)=f(x)+\frac{f^{(1)}(x)}{1!} h+\cdots+\frac{f^{(r)}(x)}{r!} h^{r}+o\left(h^{r}\right)
$$

where $o\left(h^{r}\right)$ means some (unspecified) function of $x$ and $h$ which satisfies $o\left(h^{r}\right) / h^{r} \rightarrow 0$ as $h \rightarrow 0$ (little oh notation). If $f(x)$ is $C^{r+1}$ then applying Taylor's formula at order $r+1$

$$
\begin{aligned}
f(x+h) & =f(x)+\frac{f^{(1)}(x)}{1!} h^{r}+\cdots+\frac{f^{(r)}(x)}{r!} h^{r}+\left(\frac{f^{(r+1)}(x)}{r!}+R_{r+1}(x, h)\right) h^{r+1} \\
& =f(x)+\frac{f^{(1)}(x)}{1!} h+\cdots+\frac{f^{(r)}(x)}{r!} h^{r}+O\left(h^{r+1}\right)
\end{aligned}
$$

where $O\left(h^{r+1}\right)$ denotes a continuous function times $h^{r+1}$ (big oh notation). The big oh notation has the advantage that it specifically incorporates the power $h^{r+1}$, but using it requires one more degree of smoothness than the little oh notation. There are more-or-less obvious manipulations: for example,

$$
\begin{array}{lll}
h^{2} o\left(h^{3}\right)=o\left(h^{5}\right) & \text { because } & h^{2} o\left(h^{3}\right) / h^{5}=o\left(h^{3}\right) / h^{3} \rightarrow 0 \\
o\left(h^{3}\right)^{2}=o\left(h^{6}\right) & \text { because } & o\left(h^{3}\right)^{2} / h^{6}=\left(o\left(h^{3}\right) / h^{3}\right)^{2} \rightarrow 0, \\
o\left(h^{3}\right)+o\left(h^{5}\right)=o\left(h^{3}\right) & \text { because } & \left(o\left(h^{3}\right)+o\left(h^{5}\right)\right) / h^{3}=o\left(h^{3}\right) / h^{3}+h^{2} o\left(h^{5}\right) / h^{5} \rightarrow 0
\end{array}
$$

Observe, however, that $o\left(h^{r}\right)-o\left(h^{r}\right)=0$ is false because the two little oh's generally will represent different functions. §83. The remainder in Taylor's theorem is essential in understanding the local error in polynomical approximations, This is critical where information about a function is to be extracted from its derivatives, such as where a positive Hessian at a critical point implies that the actual function has a local minimum there. The remainder may also be important when locally best-fitting two curves. For example, in elementary calculus, the tangent line is the best linear approximation to a curve at a point because it is the approximation where the error falls faster than linear. The error of the tangent approximation is $o(h)$, whereas any other linear approximation has error $o\left(h^{0}\right)$. The ratio of the first kind of error to the second kind of error falls to zero for small $h$, so the tangent line approximation is infinitely better for small $h$ than any other linear approximation.
§84. Taylor's formula is proved by iterated integration by parts as is found in single and multivariate calculus textbooks.
§85. Riemann integral: Let $f$ be a function defined on $[a, b]$. A Riemann sum is a sum of the form $\sum_{i=1}^{n} f\left(x_{i}^{*}\right) \Delta x_{i}$ where $x_{i} \in[a, b], i=0, \ldots, n$ are such that $a=x_{0}<x_{1}<x_{2}<\ldots<x_{n}=b, x_{i}^{*} \in\left[x_{i-1}, x_{i}\right]$, and $\Delta x_{i} \equiv x_{i}-x_{i-1}$. $f$ is Riemann integrable if the limit of its Riemann sums exists, in the sense that they are arbitrarily near some finite number if all the $\Delta x_{i}$ are sufficiently small, and then the limit is the integral of $f$. The principal result is that any continuous function is Riemann integrable. The following approximation result is useful:
$\S 86$. Lemma. If $e(x, u)$ is continuous, defined on an open set containing $[a, b] \times\{0\}$, and $e(x, 0)=0, x \in[a, b]$, then $\lim \sum_{i=1}^{n} e\left(x_{i}, \Delta x_{i}\right) \Delta x_{i}=0$, where the limit is in the sense of the Riemann integral.
§87. The Riemann integral is a basic tool. The approximation lemma is convenient for deriving the arc-length formula for curves from first principles and as a Riemann sum.
§88. In passing, Lemma 86 is useful in the many situations where one wants to recognize a sum of a local physical quantities as an integral, and one can approximate,

$$
\text { local physical quantity } \approx f(x) \Delta x
$$

or more precisely
local physical quantity $=f(x) \Delta x+($ error depending on $x, \Delta x)$.
For example, by definition, mass is linear density times length, in the case that the linear density is constant. In the presence of a variable linear density $\varrho(x)$, the mass of a length from $x$ to $x+\Delta x$ is $\varrho(x)$ times $\Delta x$ for small $\Delta x$, plus some small error for that accounts for the variation of the density over $\Delta x$, and the total mass is the sum of such. If the error is $e(x, \Delta x) \Delta x$, then, by Lemma 86 ,
total physical quantity $=\lim \sum_{i}\left(f\left(x_{i}\right) \Delta x_{i}+e\left(x_{i}, \Delta x_{i}\right) \Delta x\right)=\lim \sum_{i} f\left(x_{i}\right) \Delta x_{i}+\lim \sum_{i} e\left(x_{i}, \Delta x_{i}\right) \Delta x=\int_{a}^{b} f(x) d x$,
where the limits are the sense of Riemann sums, and so the total physical quantity is recognized as an integral. This same argument proves the fundamental theorem of calculus: if $f(x)$ is $C^{1}$ then, by the definition of the derivative, $f(x+u)-f(x)=f^{\prime}(x) u+e(x, u)$ where $\lim _{u \rightarrow 0} e(x, u) / u=0$, and so

$$
\begin{aligned}
f(b)-f(a) & =\lim _{x \rightarrow b^{-}} f(x)-f(a) \\
& =\lim \sum_{i} f\left(x_{i}\right)-f\left(x_{i-1}\right) \\
& =\lim \sum_{i}\left(f^{\prime}(x) \Delta x_{i}+e\left(x_{i-1}, \Delta x_{i}\right) \Delta x_{i}\right) \\
& =\lim \sum_{i} f^{\prime}(x) \Delta x_{i}+\lim \sum_{i} e\left(x_{i-1}, \Delta x_{i}\right) \Delta x_{i} \\
& =\int_{a}^{b} f^{\prime}(x) d x
\end{aligned}
$$

§89. Expositions of the Riemann integral can be found in almost any first year Calculus text. The approximation result relies on uniform continuity of a continuous function on a closed and bounded set: Given any $\epsilon>0$, uniform continuity of $e(x, u)$ on a set of the form $[a, b] \times[0, c]$ implies a $\delta$ such that $|e(x, u)|<\epsilon$ for all $x \in[a, b]$ and all $0 \leq u<\delta$. For $\max _{i} \Delta x_{i}<\delta$,

$$
\left|\sum_{i=1}^{n} e\left(t_{i}, \Delta x_{i}\right) \Delta x_{i}\right| \leq \sum_{i=1}^{n} \epsilon \Delta x_{i}=(b-a) \epsilon
$$

which is arbitrarily small for sufficiently small $\epsilon$.

## 90. The three pillars of differential geometry

$\S 91$. Theorem (Inverse function theorem). Let $f: A \subseteq \mathbb{R}^{n} \rightarrow \mathbb{R}^{n}$ be $C^{r}, r \geq 1$, let $a \in A$ and $b=f(a)$, and suppose that the matrix $D f(a)$ is nonsingular. Then there are open neighbourhoods $U \ni a$ and $V \ni b$ such that $f \mid U$ is a $C^{r}$ diffeomorphism between $U$ and $V$. Moreover, for all $x \in U, D(f \mid U)^{-1}(f(x))=(D f(x))^{-1}$.
§92. For example, the map defined by the equations

$$
u=x^{2}+\sin (x+y), \quad v=y^{3}+\tan (x y)
$$

is $C^{\infty}$ from $A \equiv\{(x, y) \mid x y \neq(2 n+1) \pi / 2\}$ to $\mathbb{R}^{2}=\{(u, v)\}$, since (1) $A$ is open because it is in the inverse image by the continuous $(x, y) \mapsto x y$ of the open set $\cdots\left(-\frac{5 \pi}{2},-\frac{3 \pi}{2}\right) \cup\left(-\frac{3 \pi}{2},-\frac{\pi}{2}\right) \cup\left(\frac{\pi}{2}, \frac{3 \pi}{2}\right) \cdots$, and (2) it is defined by a formula involving algebraic operation and trigonometric functions. But, could this map be a diffeomorphism? To analyze this directly one would have the difficult task of solving the equations for $u$ and $v$ in terms of $x$ and $y$, arriving at a formula also involving algebraic operations and trigonometric functions. The inverse function theorem provides a partial, local, answer: a map is a diffeomorphism near to some point in its domain if its derivative is invertible at that point.
§93. Given $f: A \subset \mathbb{R}^{m} \rightarrow \mathbb{R}^{m}$ and $a \in A$, choose $y$ near $a \equiv f(b)$ and define $g(x) \equiv f(x)-y$, so that inverting $f$ at $y$ is the same as solving the equation $g(x)=0$. The existence part of the proof of the inverse function theorem is obtained by showing the convergence of the (multidimensional) Newton method $x_{i+1}=x_{i}-D g(a)^{-1} g\left(x_{i}\right)$, with start $x_{1}=a$. The standard proof that the inverse has the same differentiability as the original function is more difficult, and relies on estimates of changes in the inverse directly from its definition via the Newton iteration, inductively on $r$ ([5], [6]).
§94. Vector fields and differential equations: A (time dependent) (parameterized) vector field is a map

$$
X: A \subseteq \mathbb{R}^{n} \times \mathbb{R} \times \mathbb{R}^{m}=\{(x, t, \lambda)\} \rightarrow \mathbb{R}^{n}
$$

If $(x, t, \lambda) \in A$ then $x$ is though of as a place, $t$ as a time, $\lambda$ as a parameter, and the the values of $X(x, t, \lambda)$ as a velocity; the vector field is a time and parameter dependent assignment of a velocity to place. An integral curve at parameter $\lambda$ of $X$ starting at $x$ at time $s$ is a differentiable curve function $c: I \rightarrow U$, where $I$ is an open interval, such that

$$
\begin{equation*}
\frac{d c}{d t}(t)=X(t, c(t), \lambda), \quad c(s)=x \tag{95}
\end{equation*}
$$

In words, the velocity of an integral curve must agree with that specified by the vector field. Equation (95) is the generic system of $n$ differential equations, and the term integral curve is synonymous with the solution of the initial value problem. A maximal integral curve is an integral curve $c: I \rightarrow \mathbb{R}_{\tilde{\sim}}{ }^{n}$ which does not admit an extension, i.e., $c$ is maximal if $\tilde{I}=I$ whenever $\tilde{c}: \tilde{I} \rightarrow \mathbb{R}^{n}$ is an integral curve with $I \subset \tilde{I}$ and $\tilde{c} \mid I=c$.
$\S 96$. Theorem (Ordinary differential equations existence and uniqueness). Suppose $U \subseteq \mathbb{R}^{n}$ is open, $I \subseteq \mathbb{R}$ is an open interval, $X(t, x, \lambda),(t, x, \lambda) \in I \times U \times V$, is continuous in $(t, x)$ (for each fixed $\lambda$ ), and the partial derivatives $\partial X / \partial x_{i}$ exist and are continuous. Then for every $(x, s, \lambda) \in U \times I \times V$ there is a unique maximal integral curve $c_{s, x, \lambda}: I_{s, x, \lambda} \rightarrow \mathbb{R}^{n}$ of $X$, at parameter $\lambda$ starting at $x$ at time $s$. Moreover

1. if $X$ is $C^{r}$ then the map defined by $(t, s, x, \lambda) \rightarrow c_{s, x, \lambda}(t)$ is also $C^{r}$; and
2. if $U=\mathbb{R}^{n}$ and $X$ is linear in the variable $x$ for each fixed $\lambda$ and $t$, then $I_{(s, x, \lambda)}=(a, b)$, i.e., the maximal integral curve is defined for all the times that $X$ is defined.
$\S 97$. Vector fields provide a geometrically vivid way to work with systems of differential equations and their solutions, and so have a central role in scientific modeling. Also, there are a number of objects in differential geometry that are defined as the solutions to special differential equations. Simple differential equations can have remarkably complex solutions; a number of results in differential geometry rely on the construction of complicated maps defined as the solutions of relatively simple differential equations.
§98. Local existence and uniqueness of integral curves is equivalent to local existence and uniqueness of a system of ordinary differential equations. Existence is proved by establishing the convergence of the Picard iteration

$$
c_{i+1}(t)=x+\int_{s}^{t} X\left(u, c_{i}(u), \lambda\right) d u, \quad c_{1}(t)=x
$$

Uniqueness for times near $s$ is proved using a technical result called Gronwald's Lemma, which is an (exponentially in time) estimate of the growth of deviations of initial conditions. Existence and uniqueness of maximal integral curves follows from local existence and uniqueness and the connectedness of intervals. The standard proof of smoothness of the solution in the initial conditions and the parameter relies on an induction which calculates the derivatives as solutions of related differential equations. Expositions of the theory can be found in many texts ([2], [3], [4], [8]).
§99. Invariant sets for differential equations: Let $X: U \times(a, b) \subseteq \mathbb{R}^{n} \rightarrow \mathbb{R}^{n}$ be a time dependent vector field and let $B \subseteq \mathbb{R}^{n}$. $B$ is $X$-invariant if all integral curves of $x$ starting at any point of $B$ at any time remain in $B$ for all times, i.e., if for all integral curves $x(t), x(t) \in B$ for one $t$ implies $x(t) \in B$ for all $t$. A vector $v$ is tangent to $B$ at $b \in B$ if $v=c^{\prime}(0)$ for some curve $c(t) \in B . X$ is tangent to $B$ if, for all $x$ and $t, X(x, t)$ is tangent to $B$ at $x$. The following is quite useful in proving that a set is invariant:
$\S 100$. Theorem. Let $X(x, t)$ be $C^{r}$ vector field, $r \geq 1$. Then $B \subseteq \mathbb{R}^{n}$ is $X$-invariant if $B$ is closed and $X$ is tangent to $B$.
§101. It is common to have to show that solutions of a differential equation have some property if the initial condition has that property. For any particular differential equation and property there are usually a variety of ways to do this, but Theorem 100 provides a systematic and geometrically vivid method when it is possible to define $B$ as the set of $(x, t)$ with the property.
§102. The proof uses the same Gronwald's Lemma that is used to show uniqueness of integral curves. See [1].
$\S 103$. Given functions $f_{\alpha i}\left(x_{1}, \ldots, x_{m}, y_{1}, \ldots, y_{n}\right)$, the system of partial differential equations

$$
\begin{equation*}
\frac{\partial y_{\alpha}}{\partial x_{i}}=f_{\alpha i}\left(x_{1}, \ldots, x_{m}, y_{1}, \ldots, y_{n}\right) \quad 1 \leq i \leq m, 1 \leq \alpha \leq n \tag{104}
\end{equation*}
$$

for functions $y_{\alpha}\left(x_{1}, \ldots x_{m}\right)$ occurs importantly in a number of contexts. In these equations, all first derivatives of some number of functions are specified as prescribed functions of their independent variables and the functions themselves. In contrast, the more complicated partial differential equation (the Laplace equation)

$$
\frac{\partial^{2} u}{\partial x^{2}}+\frac{\partial^{2} u}{\partial y^{2}}+\frac{\partial^{2} u}{\partial z^{2}}=0
$$

specifies the value of a mixture of partial derivatives. Equality of mixed partial derivatives implies a constraint on those $f_{\alpha i}$ for which solutions to (104) exist: substituting a solution and differentiating gives

$$
\frac{\partial^{2} y_{\alpha}}{\partial x_{j} \partial x_{i}}=\frac{\partial f_{\alpha i}}{\partial x_{j}}+\sum_{\beta} \frac{\partial f_{\alpha i}}{\partial y_{\beta}} \frac{\partial y_{\beta}}{\partial x_{j}}=\frac{\partial f_{\alpha i}}{\partial x_{j}}+\sum_{\beta} f_{\beta j} \frac{\partial f_{\alpha i}}{\partial y_{\beta}}
$$

so that existence of solutions implies, for all $\alpha, i$, and $j$, that

$$
\begin{equation*}
\frac{\partial f_{\alpha i}}{\partial x_{j}}+\sum_{\beta} f_{\beta j} \frac{\partial f_{\alpha i}}{\partial y_{\beta}}=\frac{\partial f_{\alpha j}}{\partial x_{i}}+\sum_{\beta} f_{\beta i} \frac{\partial f_{\alpha j}}{\partial y_{\beta}} \tag{105}
\end{equation*}
$$

The Frobenius theorem asserts that this necessary condition is also sufficient to solve the system (104).
$\S 106$. Theorem (Frobenius theorem). Let $f_{\alpha i}: A \times B=(x, y) \subseteq \mathbb{R}^{m} \times \mathbb{R}^{n}$ be $C^{r}$ functions, $r \geq 1$, and suppose conditions (105). Let $a \in A$ and $b \in B$. Then there are open neighbourhoods $U \ni a$ and $V \ni b$ such that, for all $u \in U$ and $v \in V$, there are unique $C^{r}$ functions $y_{\alpha}: U \rightarrow \mathbb{R}^{n}$ satisfying equations (104) and $y_{\alpha}(u)=v_{\alpha}$.
$\S 107$. The system of partial differential equations which arises as a specification of all derivatives is natural and it is not so surprising that it occurs. The applications are of great importance: for example, this is the theorem used to show that a spacetime is locally flat under the condition of vanishing Riemann tensor.
$\S 108$. The Frobenius theorem is a generalization of the standard problem of finding a local potential $\varphi$ for a vector field $X$; e.g., of finding $\varphi(x, y, z)$ such that

$$
\begin{equation*}
\frac{\partial \varphi}{\partial x}=P(x, y, z), \quad \frac{\partial \varphi}{\partial y}=Q(x, y, z), \quad \frac{\partial \varphi}{\partial z}=R(x, y, z) \tag{109}
\end{equation*}
$$

Both (104) and (109) are systems of partial differential equations; for the potential problem (109), the functions of the right side do not involve the potential itself, but rather only involve independent variables. A critical observation is that most potential problems do not have solutions, because equality of mixed partial derivative implies that, at any point $(x, y, z)$ at which a $C^{2}$ solution exists,

$$
\frac{\partial P}{\partial y}=\frac{\partial^{2} \varphi}{\partial y \partial x}=\frac{\partial^{2} \varphi}{\partial x \partial y}=\frac{\partial Q}{\partial x}, \quad \frac{\partial P}{\partial z}=\frac{\partial^{2} \varphi}{\partial z \partial x}=\frac{\partial^{2} \varphi}{\partial x \partial z}=\frac{\partial R}{\partial x}, \quad \frac{\partial Q}{\partial z}=\frac{\partial^{2} \varphi}{\partial z \partial x}=\frac{\partial^{2} \varphi}{\partial x \partial z}=\frac{\partial R}{\partial z}
$$

The restriction is important: for a vector field $X, \nabla \times X=0$ is necessary for $X$ to be conservative, and, if $\nabla \times X=0$ then a potential can be obtained from $X$ by a succession of simple integrations and substitutions. The proof of the Frobenius theorem is directly analogous, just successive solutions of ordinary differential equations replaces successive simple integrations.
§110. There is some terminology surrounding Frobenius systems: System (104) is integrable or exact if, for all $a \in A$ and $b \in B$, there are functions $y_{\alpha}$ defined on an open neighbourhood of $a$ such that $y_{\alpha i}(a)=b_{\alpha}$. It is closed or involutive if, for all $a \in A$ and all $b \in B$, there is a solution $y_{\alpha}$ defined on an open subset of $a$ and satisfying $y_{\alpha i}(a)=b$. Then Theorem 106 asserts that a $C^{1}$ Frobenius system (104) is involutive if and only if it is integrable. As mentioned in [9], the modern form of the theorem is written in terms of vector fields and Lie-brackets and the theorem above is also called the Frobenius theorem.

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