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Contents

\mathbf{C}	onter	${f ts}$	ii
In	trod	action	1
Ι	I 1	Principio d'Equivalenza	3
1	Geo	detiche	4
	1.1	Particelle non-relativistiche	4
		1.1.1 Equazione geodetica	4
	1.2	Particelle relativistiche	5
		1.2.1 Equazione geodetica	6
		1.2.2 Momento coniugato	7
		1.2.3 Interazioni	7
	1.3	Principio di Equivalenza	8
		1.3.1 Metrica di Kottler-Möller	8
		1.3.2 Gravitational time dilation	10
IJ	D	ifferential Geometry	12
2	Mai	ifolds	13
	2.1	Topological spaces	13
	2.2	Differentiable Manifolds	13
		2.2.1 Maps between manifolds	14
	2.3	Tangent spaces	14
		2.3.1 Tangent vectors	15
		2.3.2 Vector fields	16
		2.3.3 Lie derivative	18
	2.4	Tensors	19
		2.4.1 Dual Spaces	19
		2.4.2 Cotangent vectors	19
		2.4.3 Tensor fields	20
		2.4.4 Operations on tensors	21
	2.5	Differential forms	22
		2.5.1 de Rham cohomology	23

Contents

		2.5.2	Integration
3	Rie	manni	an Geometry 26
	3.1	Metric	c manifolds
		3.1.1	Metric properties
	3.2	Conne	ections
		3.2.1	Covariant derivative
		3.2.2	Covariant derivative of tensors
	3.3	Parall	el transport
		3.3.1	Normal coordinates
		3.3.2	Curvature and torsion
		3.3.3	Geodesic deviation
	3.4	Riema	ann tensor
		3.4.1	Ricci and Einstein tensors
		3.4.2	Connection and curvature forms
II	Т (Gener	ral Relativity 49
			·
4	Eins		Field Equations 50
	4.1	Einste	ein-Hilbert action
		4.1.1	Equations of motion
		4.1.2	T-0
	1.0		Diffeomorphisms
	4.2	Simple	Diffeomorphisms
	4.2	Simple 4.2.1	1
	4.2	-	e solutions
	4.2	4.2.1	e solutions
	4.2	4.2.1 4.2.2 4.2.3	e solutions
		4.2.1 4.2.2 4.2.3	e solutions 53 Minkowski space 53 de Sitter space 54 Anti-de Sitter space 57

Introduction

Nello studio delle interazioni a distanza si introducono le cosiddette teorie di campo: un campo è un'entità fisica che esiste in ogni punto dello spaziotempo (es: campo elettrico, magnetico, etc...) e che viene modificata dalla presenza di portatori della carica associata al campo.

Nel caso del di una teoria di campo per descrivere la gravità, è necessario un campo gravitazionale che sia influenzato dalla massa. Nel caso Newtoniano il campo gravitazionale $\Phi(\mathbf{r},t)$ è legato alla densità di massa $\rho(\mathbf{r},t)$ da un'equazione di Poisson:

$$\Delta \Phi = 4\pi G \rho \tag{1}$$

dove $G \approx 6.67 \cdot 10^{-11} \,\mathrm{m^3 kg^{-1} s^{-2}}$ è la costante universale di Newton.

È banale ricavare il campo gravitazionale di una massa puntiforme M: in questo caso $\rho(\mathbf{r}) = M\delta^3(\mathbf{r})$, dunque $\Phi(\mathbf{r}) = -\frac{GM}{r}$. Il caso in cui invece ρ dipende dal tempo è non banale e per essere trattato necessita di un'equazione più generale di Eq. 1: l'equazione di campo di Einstein.

Analogie e differenze con l'Elettromagnetismo Superficialmente, il problema della generalizzazione relativistica della gravitazione potrebbe sembrare analogo a quello dell'elettromagnetismo: entrambe le forze, nel caso stazionario, sono governate da una legge proporzionale a r^{-2} ed entrambi i campi sono determinati da equazioni di Poisson in cui cambia solo la costante dimensionale.

La differenza tra le due teorie di campo, però, sta proprio nella descrizione matematica delle sorgenti che subentrano nelle equazioni di Poisson: nel caso dell'Elettromagnetismo, in regime stazionario il campo elettromagnetico è determinato dalla densità di carica ρ_e e dalla densità di corrente \mathbf{J} , e per avere una descrizione relativistica bisogna combinarle in una densità di corrente quadrivettoriale $j^{\mu} = (c\rho_e, \mathbf{J})$ (si può vedere che ρ_e trasforma come una componente temporale poiché $\rho_e \sim \text{Vol}_3^{-1} \sim (\text{Vol}_4/ct)^{-1} \sim ct$, dato che il quadrivolume è un invariante di Lorentz): ciò risulta naturalmente in un potenziale quadrivettoriale $A^{\mu} = (\phi/c, \mathbf{A})$.

D'altra parte, per quanto riguarda la gravitazione, bisogna ricordare l'uguaglianza relativistica tra massa energia; inoltre, a differenza della carica elettrica, l'energia non è un invariante relativistico, ma è la componente temporale del quadrivettore impulso: in particolare, a generare il campo gravitazionale sono la densità di energia ρ e la densità di momento ρ_p , alle quali sono associate una densità di corrente di energia \mathbf{j} ed una densità di corrente di momento \mathbf{T}^i per ciascuna componente ρ_p^i . Risulta evidente che l'equazione relativistica che descrive il campo gravitazionale sia notevolmente più complicata di quella del campo elettromagnetico, poiché le sorgenti non sono descritte da un quadrivettore, bensì da un tensore, il tensore energia-impulso:

$$T^{\mu\nu} \sim \begin{bmatrix} \rho c & \rho_p c \\ \mathbf{j} & T \end{bmatrix} \tag{2}$$

Naturalmente, anche il potenziale gravitazionale sarà un tensore $h_{\mu\nu}$, ed il potenziale Newtoniano sarà $h_{00} \sim \Phi$.

2 Introduction

Scala della Relatività Generale Tramite le costanti fondamentali G e c è possibile associare ad una massa M una sua lunghezza caratteristica, detta raggio di Schwarzschild:

$$R_s := \frac{2GM}{c^2} \tag{3}$$

Le correzioni relativistiche alla teoria della gravitazione sono determinate dal parametro R_s/r e, nella maggior parte delle situazioni, sono trascurabili: basti calcolare che per la Terra $R_s \approx 10^{-2}$ m, mentre il suo raggio è $R_T = 6 \cdot 10^6$ m, dunque sulla superficie terrestre le correzioni relativistiche alla gravità Newtoniana sono dell'ordine di 10^{-8} .

Gli effetti relativistici diventano importanti quando si considerano oggetti compatti come stelle di neutroni e buchi neri.

Part I Il Principio d'Equivalenza

Geodetiche

Nelle teorie classiche di campo vengono considerati due oggetti distinti: le particelle e i campi. I campi determinano il moto delle particelle, mentre le particelle determinano le oscillazioni dei campi.

1.1 Particelle non-relativistiche

Per descrivere il moto di una particella tra due punti fissati $\mathbf{x}(t_1) \equiv \mathbf{x}_1$ e $\mathbf{x}(t_2) \equiv \mathbf{x}_2$ si studia l'azione S associata alla traiettoria $\mathbf{x}(t)$, definita come:

$$S\left[\mathbf{x}(t)\right] := \int_{t_1}^{t_2} dt \, L(\mathbf{x}(t), \dot{\mathbf{x}}(t)) \tag{1.1}$$

dove L è la lagragiana che descrive la particella.

La traiettoria percorsa dalla particella, per il principio di minima azione, è quella che estremizza S, ovvero tale per cui $\delta S = 0 \ \forall \delta \mathbf{x}(t) : \delta \mathbf{x}(t_1) = \delta \mathbf{x}(t_2) = 0$; esplicitando:

$$\delta S = \int_{t_1}^{t_2} dt \, \delta L(\mathbf{x}(t), \dot{\mathbf{x}}(t)) = \int_{t_1}^{t_2} dt \, \left(\frac{\partial L}{\partial x^i} \delta x^i + \frac{\partial L}{\partial \dot{x}^i} \delta \dot{x}^i \right)$$

$$= \int_{t_1}^{t_2} dt \, \left(\frac{\partial L}{\partial x^i} - \frac{d}{dt} \left(\frac{\partial L}{\partial \dot{x}^i} \right) \right) \delta x^i + \left[\frac{\partial L}{\partial \dot{x}^i} \delta x^i \right]_{t_1}^{t_2}$$

$$(1.2)$$

dove si è usata la convenzione di somma di Einstein.

Si vede subito che il termine di bordo è nullo, dunque estremizzare l'azione equivale alle equazioni di Eulero-Lagrange:

$$\frac{\partial L}{\partial x^i} - \frac{d}{dt} \frac{\partial L}{\partial \dot{x}^i} = 0 \tag{1.3}$$

1.1.1 Equazione geodetica

In generale, il moto di una particella libera su una generica varietà differenziale è descritto dalla lagrangiana $L = \frac{1}{2}m\dot{\mathbf{x}}\cdot\dot{\mathbf{x}}$; bisogna dunque tener conto della metrica della varietà considerata:

$$L = \frac{m}{2}g_{ij}(x)\dot{x}^i\dot{x}^j \tag{1.4}$$

dove x rappresenta collettivamente tutte le coordinate x^i sulla varietà. Si ricordi che, per una varietà reale n-dimensionale, $g_{ij} \in \mathbb{R}^{n \times n}$ è una matrice reale simmetrica.

Le equazioni di Eulero-Lagrange diventano dunque:

$$\frac{m}{2} \frac{\partial g_{ij}}{\partial x^k} \dot{x}^i \dot{x}^j - \frac{d}{dt} \left(m g_{ik} \dot{x}^i \right) = 0 \tag{1.5}$$

Espandendo il secondo termine:

$$\frac{1}{2} \frac{\partial g_{ij}}{\partial x^k} \dot{x}^i \dot{x}^j - \frac{\partial g_{ik}}{\partial x^j} \dot{x}^i \dot{x}^j - g_{ik} \ddot{x}^i = 0$$
(1.6)

Del termine $g_{ik,j} - \frac{1}{2}g_{ij,k}$, essendo contratto con un fattore simmetrico $\dot{x}^i\dot{x}^j$, sopravvive solo la parte simmetrica rispetto agli indici $i \in j$, ovvero:

$$g_{ik}\ddot{x}^i + \frac{1}{2} \left(\frac{\partial g_{ik}}{\partial x^j} + \frac{\partial g_{jk}}{\partial x^i} - \frac{\partial g_{ij}}{\partial x^k} \right) \dot{x}^i \dot{x}^j = 0$$
 (1.7)

A questo punto, si contrae per la metrica inversa g^{lk} , che per definizione soddisfa $g^{lk}g_{ik} = \delta_i^l$, così da ottenere (rinominando gli indici):

$$\ddot{x}^i + \Gamma^i_{ik}\dot{x}^j\dot{x}^k = 0 \tag{1.8}$$

dove è stato definito il simbolo di Christoffel:

$$\Gamma_{jk}^{i} := \frac{1}{2} g^{il} \left(\frac{\partial g_{lj}}{\partial x^{k}} + \frac{\partial g_{lk}}{\partial x^{j}} - \frac{\partial g_{jk}}{\partial x^{l}} \right)$$
(1.9)

Questa equazione del moto è nota come equazione geodetica e le sue soluzioni sono dette geodetiche.

1.2 Particelle relativistiche

È possibile estendere la meccanica lagrangiana allo spaziotempo di Minkowski $\mathbb{R}^{1,3}$, descritto dalla metrica:

$$\eta_{\mu\nu} = \operatorname{diag}(-1, +1, +1, +1)$$
(1.10)

Dato che questa metrica non è definita positiva, è possibile classificare due punti deparati da una distanza infinitesima $ds^2 = \eta_{\mu\nu} dx^{\mu} dx^{\nu}$ in base al segno di ds^2 : se $ds^2 < 0$ si dicono timelike-separated, se $ds^2 > 0$ spacelike-separated e se $ds^2 = 0$ lightlike-separated (o null).

A differenza del caso classico, in cui l'orbita è parametrizzata dal tempo (che è assoluto), nel caso relativistico essa deve essere parametrizzata da un generico $\sigma \in \mathbb{R}$ monotono crescente lungo la traiettoria.

In ambito relativistico, il principio di minima azione ha un'interpretazione geometrica: la traiettoria deve estremizzare la distanza tra due punti dello spaziotempo. Di conseguenza, dato che una particella di massa m deve seguire una traiettoria timelike, si definisce l'azione come:

$$S = -mc \int_{x_1}^{x_2} \sqrt{-ds^2} = -mc \int_{\sigma_1}^{\sigma_2} \sqrt{-\eta_{\mu\nu} \frac{dx^{\mu}}{d\sigma} \frac{dx^{\nu}}{d\sigma}}$$

$$\tag{1.11}$$

Il coefficiente è necessario per rendere l'azione dimensionalmente omogenea con \hbar . L'azione così definita presenta due simemtrie:

1. invarianza di Lorentz: l'azione è invariante per $x^{\mu} \mapsto \Lambda^{\mu}_{\nu} x^{\nu}$, con $\Lambda : \Lambda^{\mu}_{\rho} \eta_{\mu\nu} \Lambda^{\nu}_{\sigma} = \eta_{\rho\sigma}$;

2. invarianza per riparametrizzazioni: essendo σ un parametro arbitrario, è normale che l'azione non dipenda dalla sua scelta, infatti se si riparametrizza con una funzione monotona $\tilde{\sigma}(\sigma)$ si ha:

$$\tilde{S} = -mc \int_{\tilde{\sigma}_1}^{\tilde{\sigma}_2} d\tilde{\sigma} \sqrt{-\eta_{\mu\nu} \frac{dx^{\mu}}{d\tilde{\sigma}} \frac{dx^{\nu}}{d\tilde{\sigma}}} = -mc \int_{\sigma_1}^{\sigma_2} d\sigma \frac{d\tilde{\sigma}}{d\sigma} \sqrt{-\eta_{\mu\nu} \frac{dx^{\mu}}{d\sigma} \frac{dx^{\nu}}{d\sigma} \left(\frac{d\sigma}{d\tilde{\sigma}}\right)^2} = S \qquad (1.12)$$

Grazie all'invarianza per riparametrizzazioni, il valore dell'azione tra due punti dello spaziotempo assume un significato ben preciso, il tempo proprio, ovvero il tempo misurato dalla particella in moto stessa:

$$\tau(\sigma) = \frac{1}{c} \int_0^{\sigma} d\sigma' \sqrt{-\eta_{\mu\nu} \frac{dx^{\mu}}{d\sigma'}} \frac{dx^{\nu}}{d\sigma'}$$
(1.13)

Una conseguenza dell'identificazione tra azione e tempo proprio è che il principio di minima azione richiede che la traiettoria estremizzi il tempo proprio. È anche possibile riparametrizzare l'azione col tempo proprio, essendo questo una funzione monotona crescente lungo la traiettoria.

1.2.1 Equazione geodetica

Nel caso relativistico su una varietà differenziabile generica, la lagrangiana di una particella libera è:

$$L = \sqrt{-g_{\mu\nu}(x)\dot{x}^{\mu}\dot{x}^{\nu}} \tag{1.14}$$

Dunque, le equazioni di Eulero-Lagrange diventano:

$$-\frac{1}{2L}\frac{\partial g_{\mu\nu}}{\partial x^{\rho}}\dot{x}^{\mu}\dot{x}^{\nu} - \frac{d}{d\sigma}\left(-\frac{1}{L}g_{\rho\nu}\dot{x}^{\nu}\right) = 0 \tag{1.15}$$

L'unica differenza con Eq. 1.5 è che $L=L(\sigma)$, dunque si trova un'equazione analoga all'Eq. 1.7 ma con un termine aggiuntivo:

$$g_{\mu\rho}\ddot{x}^{\rho} + \frac{1}{2} \left(\frac{\partial g_{\mu\rho}}{\partial x^{\nu}} + \frac{\partial g_{\mu\nu}}{\partial x^{\rho}} - \frac{\partial g_{\nu\rho}}{\partial x^{\mu}} \right) \dot{x}^{\nu} \dot{x}^{\rho} = \frac{1}{L} \frac{dL}{d\sigma} g_{\mu\rho} \dot{x}^{\rho}$$
(1.16)

È possibile annullare il termine $\frac{dL}{d\sigma}$ con un'opportuna scelta di parametrizzazione. Dall'Eq. 1.13 si vede che:

$$c\frac{d\tau}{d\sigma} = L(\sigma) \tag{1.17}$$

Dunque, riparametrizzando con $\tau(\sigma)$:

$$L(\tau) = \sqrt{-g_{\mu\nu}(x)\frac{dx^{\mu}}{d\tau}\frac{dx^{\nu}}{d\tau}} = \frac{d\sigma}{d\tau}L(\sigma) = c$$
 (1.18)

In generale, qualsiasi riparametrizzazione con $\tilde{\tau} = a\tau + b$ (parametri affini della worldline) porta ad avere una lagrangiana costante.

Ricordando la definizione di connessione affine in Eq. 1.9, si trova l'equazione geodetica:

$$\frac{d^2x^{\mu}}{d\tau^2} + \Gamma^{\mu}_{\nu\rho} \frac{dx^{\nu}}{d\tau} \frac{dx^{\rho}}{d\tau} \tag{1.19}$$

1.2.2 Momento coniugato

La differenza sostanziale tra lo spazio euclideo e lo spaziotempo di Minkowski è che, mentre nello spazio euclideo un corpo può rimanere fermo, nello spaziotempo nessun corpo può fermarsi nella direzione temporale. Questo fatto deve essere rispecchiato dal momento della particella:

$$p_{\mu} = \frac{dL}{dx^{\mu}} = \frac{d}{dx^{\mu}} \left(-mc\sqrt{-\eta_{\mu\nu}} \dot{x}^{\mu} \dot{x}^{\nu} \right) = mc \frac{\eta_{\mu\nu} \dot{x}^{\nu}}{\sqrt{-\eta_{\rho\sigma}} \dot{x}^{\rho} \dot{x}^{\sigma}} = -\frac{m^{2}c^{2}}{L} \eta_{\mu\nu} \dot{x}^{\nu}$$
(1.20)

Non tutte le componenti del 4-momento sono indipendenti:

$$p^{2} = p^{\mu}p_{\mu} = \frac{m^{4}c^{4}}{L^{2}}\eta_{\mu\nu}\dot{x}^{\mu}\dot{x}^{\nu} = -m^{2}c^{2}$$
(1.21)

$$(p^0)^2 = \mathbf{p}^2 + m^2 c^2 \tag{1.22}$$

Di conseguenza, si ha sempre $(p^0)^2 > 0$.

Si noti che riparametrizzando la worldline col tempo proprio, dato che $\frac{d\tau}{d\sigma} = -\frac{L}{mc^2}$:

$$p^{\mu} = m \frac{d\sigma}{d\tau} \frac{dx^{\mu}}{d\sigma} = m \frac{dx^{\mu}}{d\tau} \tag{1.23}$$

La non-indipendenza di una delle componenti del 4-momento è naturale: da una descrizione classica del sistema risultano tre gradi di libertà $x^i(t)$, dunque passando ad una descrizione relativistica non può risultare un grado di libertà in più. Ciò è legato all'invarianza per riparametrizzazione: risolvendo le equazioni del moto si trovano le componenti della traiettoria $x^{\mu} = x^{\mu}(\sigma)$, ma il parametro σ non può rappresentare dell'informazione sul sistema, dunque una delle quattro equazioni del moto va utilizzata per eliminare la dipendenza da σ , riducendo di nuovo a tre i gradi di libertà.

1.2.3 Interazioni

Dall'invarianza per riparametrizzazione, è possibile scegliere come parametro $\sigma=t$ il tempo misurato in un qualunque RF inerziale; considerando una particella libera nello spaziotempo di Minkowski, l'azione in Eq. 1.11 diventa:

$$S = -mc^2 \int_{t_0}^{t_1} dt \sqrt{1 - \frac{\dot{\mathbf{x}}^2}{c^2}}$$
 (1.24)

In questa forma, è chiara la presenza di soli tre gradi di libertà dovuti a $\mathbf{x}(t)$.

1.2.3.1 Elettromagnetismo

Analogamente al caso classico, per rappresentare l'interazione elettromagnetica è necessario aggiungere un termine potenziale all'azione. Il problema è che la semplice aggiunta di $\int d\sigma V(\mathbf{x})$ non soddisfa l'invarianza per riparametrizzazione: per soddisfarre questo requisito, è necessario individuare un potenziale che cancelli il fattore Jacobiano derivante dalla trasformazione della misura $d\sigma$. Un'opzione è considerare un termine lineare in \dot{x}^{μ} , dunque per l'invarianza di Lorentz è necessario che l'indice μ sia contratto:

$$S = \int_{\sigma_1}^{\sigma_2} d\sigma \left[-mc^2 \sqrt{-\eta_{\mu\nu}} \frac{dx^{\mu}}{d\sigma} \frac{x^{\nu}}{d\sigma} - qA_{\mu}(x)\dot{x}^{\mu} \right]$$
 (1.25)

dove q è la carica associata all'interazione e $A_{\mu}(x)$ è il suo potenziale quadrivettoriale. Scrivendo $A_{\mu}(x) = (\phi(x)/c, \mathbf{A}(x))$, l'azione in Eq. 1.25 descrive l'interazione elettromagnetica; ciò diventa evidente riparametrizzando con $\sigma = t$:

$$S = \int_{t_0}^{t_1} dt \left[-mc^2 \sqrt{1 - \frac{\dot{\mathbf{x}}^2}{c^2}} - q\phi(x) - q\mathbf{A}(x) \cdot \mathbf{x} \right]$$
 (1.26)

1.2.3.2 Gravitazione

Per descrivere l'interazione gravitazione è necessario considerare un'azione generalizzata del tipo:

$$S = \int_{t_0}^{t_1} dt \left[-mc^2 \sqrt{1 + \frac{2\Phi(\mathbf{x})}{c^2} - \frac{\dot{\mathbf{x}}^2}{c^2}} \right]$$
 (1.27)

Nel limite non-relativistico $\dot{\mathbf{x}}^2 \ll c^2$ e $2\Phi(\mathbf{x}) \ll c^2$, dunque approssimando al prim'ordine:

$$S = \int_{t_0}^{t_1} dt \left[-mc^2 + \frac{m}{2}\dot{\mathbf{x}}^2 - m\Phi(\mathbf{x}) \right]$$
 (1.28)

Il primo termine (l'energia a riposo della particella) non ha effetti sull'azione poiché è costante, mentre gli altri termini descrivono il moto non-relativistico di una particella in un campo gravitazionale $\Phi(\mathbf{x})$. Il termine $1 + 2\Phi(\mathbf{x})/c^2$ in Eq. 1.27 deriva dalla componente η_{00} della metrica, dunque si osserva che la metrica deve dipendere da x, ovvero la descrizione dell'interazione gravitazionale introduce uno spaziotempo curvo. La condizione che deve soddisfarre la metrica in un weak gravitational field è:

$$g_{00}(x) \approx -\left(1 + \frac{2\Phi(x)}{c^2}\right) \tag{1.29}$$

con $\Phi(x)$ il campo gravitazionale Newtoniano.

1.3 Principio di Equivalenza

Come si evince dall'Eq. 1.28, la "carica" dell'interazione gravitazionale è porprio la massa della particella: questo fatto viene definito weak equivalence principle (WEP) ed è solitamente espresso tramite l'uguaglianza tra la massa inerziale e la massa gravitazionale:

$$m_{\rm i} = m_{\rm g} \tag{1.30}$$

Questo è un fatto sperimentale misurato con una precisione dell'ordine di 10^{-13} .

1.3.1 Metrica di Kottler-Möller

Una conseguenza del WEP è l'indistinguibilità tra un'accellerazione costante ed un campo gravitazionale costante: ciò può essere visto in maniera analitica.

Considerando una particella di massa m con accellerazione costante $\mathbf{a} = a\hat{\mathbf{e}}_x$ in un RF inerziale \mathcal{O} , dalla relatività speciale si vede subito che la traiettoria non può essere $x(t) = \frac{1}{2}at^2$, poiché la velocità eccederebbe c; bisogna invece ricordare la composizione relativistica delle velocità:

$$v = \frac{v_1 + v_2}{1 + v_1 v_2 / c^2} \tag{1.31}$$

È possibile definire la rapidità $\varphi : v = c \tanh \varphi$, così da poter riscrivere la composizione delle velocità come $\varphi = \varphi_1 + \varphi_2$.

Un'accelerazione costante significa che la rapidità della particella aumenta linearmente rispetto al tempo proprio, ovvero $\varphi(\tau) = a\tau/c$, quindi:

$$v(\tau) = \frac{dx}{d\tau} = c \operatorname{sech}\left(\frac{a\tau}{c}\right) \tag{1.32}$$

La relazione tra il tempo misurato nel RF inerziale ed il tempo proprio è:

$$\frac{dt}{d\tau} = \gamma(\tau) = \sqrt{\frac{1}{1 - v(\tau)^2/c^2}} = \cosh\left(\frac{a\tau}{c}\right) \implies t(\tau) = \frac{c}{a}\sinh\left(\frac{a\tau}{c}\right) \tag{1.33}$$

con costante d'integrazione tale per cui $\tau=0$ corrisponda a t=0. Per quanto riguarda la traiettoria:

$$x(\tau) = \frac{c^2}{a} \cosh\left(\frac{a\tau}{c}\right) - \frac{c^2}{a} \tag{1.34}$$

con costante d'integrazione tale per cui x(0) = 0. Si trova dunque un'iperbole nello spaziotempo:

$$\left(x + \frac{c^2}{a}\right)^2 - c^2 t^2 = \frac{c^4}{a^2} \tag{1.35}$$

con asintoti $ct = \pm (x + c^2/a)$ per $\tau \to \pm \infty$. Si può mostrare che la trasformazione tra le coordinate (t,x) nel RF inerziale e quelle (τ,ρ) solidali alla particella è data da:

$$ct = \left(\rho + \frac{c^2}{a}\right) \sinh\left(\frac{a\tau}{c}\right)$$

$$x = \left(\rho + \frac{c^2}{a}\right) \cosh\left(\frac{a\tau}{c}\right) - \frac{c^2}{a}$$
(1.36)

Infatti, l'orbita della particella è giustamente descritta da $\rho = 0$. Inoltre, si vede che le coordinate (τ, ρ) non ricoprono tutto lo spaziotempo di Minkowski (t, x): questo dimostra che ci sono delle regioni dello spaziotempo causalmente disconnesse dalla particella.

Ricordando che $ds^2 = -c^2 dt^2 + d\mathbf{r}^2$, sostituendo l'Eq. 1.36 si ottiene la metrica di Kottler-Möller:

$$ds^{2} = -\left(1 + \frac{a\rho}{c^{2}}\right)^{2}c^{2}d\tau^{2} + d\rho^{2} + dy^{2} + dz^{2}$$
(1.37)

Mentre la parte spaziale rimane piatta, si vede che $g_{00} = g_{00}(x)$; inoltre, nel caso sub-relativistico:

$$g_{00} \approx -\left(1 + \frac{2a\rho}{c^2}\right) \tag{1.38}$$

Definendo $\Phi(\rho) = a\rho$, si trova proprio la condizione 1.29: questo è proprio l'assero del WEP, poiché un'accellerazione costante dà una metrica indistinguibile da quella di un campo gravitazionale costante (sub-relativistico).

1.3.1.1 Principio di equivalenza di Einstein

Dal WEP deriva il fatto che un campo gravitazionale costante può essere annullato dalla scelta di un particolare RF, il free-fall RF.

Il principio di equivalenza di Einstein è una generalizzazione del WEP: esso afferma che esiste sempre un local RF in cui gli effetti di un qualsiasi campa gravitazionali sono localmente annullati. Formalmente, ciò equivale a dire che la metrica $g_{\mu\nu}$ è sempre localmente approssimabile con la metrica di Minkowski $\eta_{\mu\nu}$.

Gli effetti di un campo gravitazionale non uniforme diventano evidenti quando è possibile svolgere misurazioni su una regione estesa di spazio. Si consideri ad esempio un osservatore confinato in un cubo chiuso in free-fall verso la Terra: non esiste alcun esperimento locale in grado di distinguere se l'osservatore stia fluttuando nello spazio oppure sia in free-fall, bensì è necessario un esperimento non-locale; un esempio di questo tipo di esperimenti consiste nel lasciare libere due masse test (quindi non influenzate vicendevolmente per interazione gravitazionale) separate da una certa distanza: se il cubo sta fluttuando nello spazio, le masse rimaranno nella loro posizione iniziale per inerzia, mentre se esso è in free-fall ciascuna massa sarà attratta verso il centro della Terra, dunque il loro spostamento avrà non solo una componente verticale (rispetto alla caduta), ma anche una orizzontale, per quanto piccola: le masse si sposteranno dunque l'una verso l'altra per effetto di una tidal force, uno dei principali fattori discriminanti dei campi gravitazionali non uniformi.

1.3.2 Gravitational time dilation

In condizioni di campo gravitazionale debole si ha $g_{00}(x) = 1 + \frac{2\Phi(x)}{f^2}$. Considerando il campo gravitazionale di un corpo sferico di massa M uniforme, $\Phi(r) = -\frac{GM}{r}$, si ha che un osservatore ad una distanza fissa r misurerà degli intervalli di tempo dati da:

$$d\tau^2 = g_{00}dt^2 = \left(1 - \frac{2GM}{rc^2}\right)dt^2 \tag{1.39}$$

Dunque, definendo t il tempo misurato da un osservatore a $r \to \infty$, l'osservatore a r misurerà:

$$T(r) = t\sqrt{1 - \frac{2GM}{rc^2}}\tag{1.40}$$

ovvero il tempo scorre più lentamente vicino ad un corpo massivo.

È anche possibile mettere in relazione i tempi misurati a due distanze finite r_1 ed $r_2 = r_1 + \Delta r$:

$$T_{2} = t\sqrt{1 - \frac{2GM}{(r_{1} + \Delta r)c^{2}}} \approx t\sqrt{1 - \frac{2GM}{r_{1}c^{2}} + \frac{2GM\Delta r}{r_{1}^{2}c^{2}}}$$

$$\approx t\sqrt{1 - \frac{2GM}{r_{1}c^{2}}} \left(1 + \frac{GM\Delta r}{r_{1}^{2}c^{2}}\right) = T_{1}\left(1 + \frac{GM\Delta r}{r_{1}^{2}c^{2}}\right)$$
(1.41)

dove si è assunto che $\Delta r \ll r_1$ e $2GM \ll r_1c^2$. Ad esempio, con un per un dislivello di 10^3 m sul livello del mare ($\sim 6 \cdot 10^6$ m) si ha una differenza di 10^{-12} s, confermato sperimentalmente.

1.3.2.1 Gravitational redshift

Un'importante conseguenza della gravitational time dilation è il gravitational redshift. Sempre in condizioni di campo debole, si consideri un segnale a distanza r_1 che si ripete ad intervalli ΔT_1 : un osservatore a r_2 misurerà:

$$\Delta T_2 = \sqrt{\frac{1 + 2\Phi(r_2)/c^2}{1 + 2\Phi(r_1)/c^2}} \Delta T_1 \approx \left(1 + \frac{\Phi(r_2) - \Phi(r_1)}{c^2}\right) \Delta T_1$$
 (1.42)

Dato che $\omega \sim T^{-1},$ si ha:

$$\omega_2 \approx \left(1 + \frac{\Phi(r_2) - \Phi(r_1)}{c^2}\right)^{-1} \omega_1 \tag{1.43}$$

Dato che $\Phi(r) \sim r^{-1}$ (nel caso considerato), se $r_2 > r_1$ si ha $\omega_2 < \omega_1$ (redshift), mentre se $r_2 < r_1$ si ha $\omega_2 > \omega_1$ (blueshift).

Part II Differential Geometry

Manifolds

2.1 Topological spaces

Definition 2.1.1. The topology \mathcal{T} of a set X is a family of subsets of X, i.e. $\mathcal{T} \subseteq \mathcal{P}(X)$, defined as open sets, with the following properties:

- 1. $\emptyset, X \in \mathcal{T}$;
- 2. $O_{\alpha}, O_{\beta} \in \mathcal{T} \Rightarrow O_{\alpha} \cap O_{\beta} \in \mathcal{T}$;
- 3. $\{O_{\alpha}\}_{{\alpha}\in I}\subset \mathcal{T}\ (I \text{ arbitrary index set}) \Rightarrow \bigcup_{{\alpha}\in I} O_{\alpha}\in \mathcal{T}.$

Definition 2.1.2. A topological space M is a set of points, endowed with a topology \mathcal{T} .

Definition 2.1.3. Given a topological space (M, \mathcal{T}) , $O \in \mathcal{T}$ is a *neighbourhood* of a point $p \in M$ if $p \in O$.

Definition 2.1.4. A topological space (M, \mathcal{T}) is Hausdorff if $\forall p, q \in M \exists O_1, O_2 \in \mathcal{T}$ neighbourhoods of p and q respectively such that $O_1 \cap O_2 = \emptyset$.

Definition 2.1.5. A homeomorphism between two topological spaces (M_1, \mathcal{T}_1) and (M_2, \mathcal{T}_2) is a bijective map $f: M_1 \to M_2$ which is bicontinuous, i.e. both f and f^{-1} are continuous: f is continuous if $O \in \mathcal{T}_2 \Rightarrow f^{-1}(O) \in \mathcal{T}_1$.

2.2 Differentiable Manifolds

Definition 2.2.1. An *n*-dimensionale differentiable manifold \mathcal{M} is a Hausdorff topological space such that:

- 1. \mathcal{M} is locally homeomorphic to \mathbb{R}^n , i.e. $\forall p \in \mathcal{M} \exists O \in \mathcal{T}(\mathcal{M}) : p \in O \land \exists \varphi : O \to U \in \mathcal{T}(\mathbb{R}^n)$ homeomorphism;
- 2. given $O_{\alpha}, O_{\beta} \in \mathcal{T}(\mathcal{M}) : O_{\alpha} \cap O_{\beta} \neq \emptyset$, the corresponding maps $\varphi_{\alpha} : O_{\alpha} \to U_{\alpha}, \varphi_{\beta} : O_{\beta} \to U_{\beta}$ must be *compatible*, i.e. $\varphi_{\beta} \circ \varphi_{\alpha}^{-1} : \varphi_{\alpha}(O_{\alpha} \cap O_{\beta}) \to \varphi_{\beta}(O_{\alpha} \cap O_{\beta})$ and its inverse must be smooth (of \mathcal{C}^{∞} class).

The maps φ_{α} are called *charts* and a collection of compatible charts is called an *atlas*: a maximal atlas \mathcal{A} is an atlas such that $\bigcup_{\alpha \in I} O_{\alpha} = \mathcal{M}$. Two atlases are compatible if each chart of one atlas is compatible with every chart of the other: they define the same differentiable structure on the manifold.

Each chart φ_{α} provides a coordinate system on the region O_{α} : $\varphi_{\alpha}(p) = (x^{1}(p), \dots, x^{\mu}(p), \dots, x^{n}(p))$. The transition functions $\varphi_{\beta} \circ \varphi_{\alpha}^{-1}$ are therefore coordinate transformations on overlapping regions.

Example 2.2.1. The n-sphere \mathbb{S}^n is a differentiable manifold.

Example 2.2.2. To define a differentiable structure on S^1 an atlas of two charts is needed: the standard parametrization $\theta \in [0, 2\pi)$ is not a well-defined chart because $[0, 2\pi)$ is not an open set in the Euclidean topology of \mathbb{R} , therefore the elimination of a point is necessary; usually, the two charts of the atlas are defined by $\theta_1 \in (0, 2\pi)$, excluding (1, 0) (in the embedding space \mathbb{R}^2), and $\theta_2 \in (-\pi, \pi)$, excluding (-1, 0): they are evidently compatible, thus they form a maximal atlas.

2.2.1 Maps between manifolds

Locally mapping \mathcal{M} to \mathbb{R}^n allows to import concepts of Analysis from \mathbb{R}^n to \mathcal{M} .

Definition 2.2.2. A function $f: \mathcal{M} \to \mathbb{R}$ on a differentiable manifold $(\mathcal{M}, \mathcal{A})$ is *smooth* if $f \circ \varphi_{\alpha}^{-1}$: $U_{\alpha} \to \mathbb{R}$ is smooth for all charts $(U_{\alpha}, \varphi_{\alpha}) \in \mathcal{A}$.

Definition 2.2.3. A map $f: \mathcal{M} \to \mathcal{N}$ between two differentiable manifolds $(\mathcal{M}, \mathcal{A}_1), (\mathcal{N}, \mathcal{A}_2)$ is smooth if $\psi_{\alpha_2} \circ f \circ \varphi_{\alpha_1}^{-1}: U_{\alpha_1} \to V_{\alpha_2}$ is smooth for all charts $(U_{\alpha_1}, \varphi_{\alpha_1}) \in \mathcal{A}_1, (V_{\alpha_2}, \varphi_{\alpha_2}) \in \mathcal{A}_2$.

Definition 2.2.4. A diffeomorphism between two differentiable manifolds \mathcal{M}, \mathcal{N} is a smooth homeomorphism $f : \mathcal{M} \to \mathcal{N}$.

Proposition 2.2.1. If \mathcal{M} and \mathcal{N} are diffeomorphic, then $\dim_{\mathbb{R}} \mathcal{M} = \dim_{\mathbb{R}} \mathcal{N}$.

Example 2.2.3. \mathbb{S}^7 can be covered by multiple incompatible at lases: the resulting manifolds are homeomorphic but not diffeomorphic.

Example 2.2.4. \mathbb{R}^n has a unique differentiable structure for all $n \in \mathbb{N}$, except for n = 4: \mathbb{R}^4 can be covered by infinitely-many incompatible atlases.

2.3 Tangent spaces

The notions of calculus can be defined on a differential manifold $(\mathcal{M}, \mathcal{A})$ via tangent spaces.

Definition 2.3.1. The derivative of a function $f : \mathcal{M} \to \mathbb{R}$ at a point $p \in \mathcal{M}$, covered by the chart (φ, U) , is defined as:

$$\left. \frac{\partial f}{\partial x^{\mu}} \right|_{p} := \left. \frac{\partial (f \circ \varphi^{-1})}{\partial x^{\mu}} \right|_{\varphi(p)} \tag{2.1}$$

Evidently, this definition depends on the choise of coordinates x^{μ} , thus it depends on the chart.

2.3.1 Tangent vectors

Definition 2.3.2. The set of all smooth functions on \mathcal{M} is denoted by $\mathcal{C}^{\infty}(\mathcal{M})$.

Definition 2.3.3. A tangent vector to \mathcal{M} in $p \in \mathcal{M}$ is an operator $X_p : \mathcal{C}^{\infty}(\mathcal{M}) \to \mathbb{R}$ such that:

- 1. $X_p(f+g) = X_p(f) + X_p(g) \forall f, g \in \mathcal{C}^{\infty}(\mathcal{M});$
- 2. $X_p(f) = 0$ for all constant functions;
- 3. $X_p(fg) = X_p(f)g(p) + f(p)X_p(g) \,\forall f, g \in \mathcal{C}^{\infty}(\mathcal{M}).$

Proposition 2.3.1. $X_p(\alpha f) = \alpha X_p(f) \, \forall \alpha \in \mathbb{R}.$

Proof. Trivial from conditions 2. and 3. of Def. 2.3.3.

It is simple to check that $\partial_{\mu}|_{p}$ satisfies the conditions of Def. 2.3.3.

Theorem 2.3.1. The set $T_p\mathcal{M}$ of all tangent vectors at a point $p \in \mathcal{M}$ forms an n-dimensional space, called tangent space, and $\{\partial_{\mu}|_p\}_{\mu=1,\dots,n}$ is a basis of such space.

Proof. Defining $f \circ \varphi^{-1} \equiv F : U \subset \mathcal{M} \to \mathbb{R}$, with $f : \mathcal{M} \to \mathcal{M}$ and $(\varphi, U) \in \mathcal{A}$, it can be proved that, in some neighbourhood of p (not necessarily U), F call always be written as:

$$F(x) = F(x^{\mu}(p)) + (x^{\mu} - x^{\mu}(p)) F_{\mu}(x)$$

for some n functions F_{μ} (ex.: a Taylor series, or more generally $F(x) = F(0) + x \int_0^1 dt \, F(xt)$). Applying $\partial_{\mu}|_{x(p)}$:

$$\left. \frac{\partial F}{\partial x^{\mu}} \right|_{x(p)} = F_{\mu}(x(p))$$

Defining $f_{\mu} \equiv F_{\mu} \circ \varphi$, for any $q \in \mathcal{M}$ in an appropriate neighbourhood of p:

$$f(q) = f(p) + (x^{\mu}(q) - x^{\mu}(p)) f_{\mu}(q)$$

Moreover, remembering Eq. 2.1:

$$f_{\mu}(p) = F_{\mu} \circ \varphi(p) = F_{\mu}(x(p)) = \frac{\partial F}{\partial x^{\mu}}\Big|_{x(p)} = \frac{\partial f}{\partial x^{\mu}}\Big|_{p}$$

Using these facts, the action of a tangent vector can be written explicitly:

$$X_{p}(f) = X_{p} (f(p) + (x^{\mu} - x^{\mu}(p)) f_{\mu})$$

$$= X_{p} (f(p)) + X_{p} ((x^{\mu} - x^{\mu}(p))) f_{\mu}(p) + (x^{\mu} - x^{\mu}(p)) (p) X_{p} (f_{\mu})$$

$$= X_{p} (x^{\mu}) f_{\mu}(p)$$

because f(p) is a constant and $(x^{\mu} - x^{\mu}(p))(p) = x^{\mu}(p) - x^{\mu}(p) = 0$. Therefore, remembering the expression for $f_{\mu}(p)$:

$$X_p = X_p(x^\mu) \frac{\partial}{\partial x^\mu} \bigg|_p \equiv X^\mu \frac{\partial}{\partial x^\mu} \bigg|_p$$

Thus, $T_p\mathcal{M} = \operatorname{span}\{\partial_{\mu}|_p\}$. To check for linear independence, suppose $\alpha = \alpha^{\mu}\partial_{\mu}|_p \equiv 0$: acting on $f = x^{\nu}$, it gives $\alpha(f) = \alpha_{\mu}\partial_{\mu}(x^{\nu})|_p = \alpha_{\nu} = 0$. This concludes the proof.

2.3.1.1 Changing coordinates

Although $\partial_{\mu}|_{p}$ depends on the choice of coordinates (it is a *coordinate basis*), the existence of X_{p} is independent of that choice.

If two different charts (φ, U) , $(\tilde{\varphi}, V)$ intersect in a neighbourhood of $p \in U \cap V$, the transition from x^{μ} to y^{μ} can be expressed as:

$$X_{p}(f) = X^{\mu} \frac{\partial f}{\partial x^{\mu}} \bigg|_{p} = X^{\mu} \frac{\partial y^{\nu}}{\partial x^{\mu}} \bigg|_{\varphi(p)} \frac{\partial f}{\partial y^{\nu}} \bigg|_{p}$$

$$(2.2)$$

This equation can have two interpretations: the alibi interpretation:

$$\frac{\partial}{\partial x^{\mu}}\bigg|_{p} = \frac{\partial y^{\nu}}{\partial x^{\mu}}\bigg|_{\varphi(p)} \frac{\partial}{\partial y^{\nu}}\bigg|_{p} \tag{2.3}$$

and the alias interpretation:

$$\tilde{X}^{\nu} = X^{\mu} \frac{\partial y^{\nu}}{\partial x^{\mu}} \bigg|_{\varphi(p)} \tag{2.4}$$

Components of vectors which transform this way are called *contravariant*.

2.3.1.2 Curves

Consider a smooth curve on \mathcal{M} , i.e. a smooth map $\sigma: I \in \mathcal{T}(\mathbb{R}) \to \mathcal{M}$, parametrized as $\sigma(t): \sigma(0) = p \in \mathcal{M}$; with a given chart (φ, U) , this curve becomes $\varphi \circ \sigma: I \to \mathbb{R}^n$, parametrized by $x^{\mu}(t)$. The tangent vector to the curve in p is:

$$X_{p} = \frac{dx^{\mu}(t)}{dt} \bigg|_{t=0} \frac{\partial}{\partial x^{\mu}} \bigg|_{p} \tag{2.5}$$

This operator, applied to a function $f \in \mathcal{C}^{\infty}(\mathcal{M})$, calculates the directional derivative of f along the curve. It can be showed that every tangent vector can be written as in Eq. 2.5, therefore the tangent space is literally the space of all possible tangents to curves passing through p.

It must be noted that tangent spaces at different points are entirely different spaces: there's no way to directly compare vectors between them.

2.3.2 Vector fields

Definition 2.3.4. A vector field X is a smooth map $X: p \in \mathcal{M} \mapsto X_p \in T_p\mathcal{M}$. It can also be viewed as a smooth map $X: \mathcal{C}^{\infty}(\mathcal{M}) \to \mathcal{C}^{\infty}(\mathcal{M})$, as $(X(f))(p) = X_p(f) \in \mathbb{R}$.

Definition 2.3.5. The space of all vector fields on \mathcal{M} is denoted by $\mathfrak{X}(\mathcal{M})$.

Given a chart (φ, U) , a vector field X can be expressed as:

$$X = X^{\mu} \frac{\partial}{\partial x^{\mu}} \tag{2.6}$$

with $X^{\mu} \in \mathcal{C}^{\infty}(\mathcal{M})$. This expression is only defined on U.

2.3.2.1 Lie brakets

Given two vector fields $X, Y \in \mathfrak{X}(\mathcal{M})$, their product is clearly not a vector field, as it does not satisfy Leibniz' rule:

$$XY(fg) = XY(f)g + Y(f)X(g) + X(f)Y(g) + fXY(g) \neq XY(f)g + fXY(g)$$

where $XY(f) \equiv X(Y(f))$.

Definition 2.3.6. Given two vector fields $X, Y \in \mathfrak{X}(\mathcal{M})$, their *commutator* (or *Lie bracket*) is defined as:

$$[X,Y](f) = XY(f) - YX(f)$$
 (2.7)

With a given chart:

$$[X,Y](f) = X^{\mu} \frac{\partial}{\partial x^{\mu}} \left(Y^{\nu} \frac{\partial f}{\partial x^{\nu}} \right) - Y^{\mu} \frac{\partial}{\partial x^{\mu}} \left(X^{\nu} \frac{\partial f}{\partial x^{\nu}} \right)$$
$$= \left(X^{\mu} \frac{\partial Y^{\nu}}{\partial x^{\mu}} - Y^{\mu} \frac{\partial X^{\nu}}{\partial x^{\mu}} \right) \frac{\partial f}{\partial x^{\nu}}$$

therefore:

$$[X,Y] = \left(X^{\mu} \frac{\partial Y^{\nu}}{\partial x^{\mu}} - Y^{\mu} \frac{\partial X^{\nu}}{\partial x^{\mu}}\right) \frac{\partial}{\partial x^{\nu}}$$
 (2.8)

Theorem 2.3.2 (Jacobi). Given $X, Y, Z \in \mathfrak{X}(\mathcal{M})$, the Jacobi identity holds:

$$[X, [Y, Z]] + [Y, [Z, X]] + [Z, [X, Y]] = 0 (2.9)$$

Proposition 2.3.2. $\mathfrak{X}(\mathcal{M})$ is a Lie algebra.

2.3.2.2 Integral curves

Definition 2.3.7. A *flow* on \mathcal{M} is a one-parameter family of diffeomorphisms $\sigma_t : \mathcal{M} \to \mathcal{M}$, labelled by $t \in \mathbb{R}$, with group structure: $\sigma_0 = \mathrm{id}_{\mathcal{M}}$ and $\sigma_s \circ \sigma_t = \sigma_{s+t}$, thus $\sigma_{-t} = \sigma_t^{-1}$.

Such flows give rise to streamlines on the manifold: these streamlines are required to be smooth. Defining $x^{\mu}(\sigma_t) \equiv x^{\mu}(t)$, a vector field can be defined by the tangent to the streamlines at each point on the manifold:

$$X^{\mu}(x^{\mu}(t)) = \frac{dx^{\mu}(t)}{dt} \tag{2.10}$$

The inverse reasoning is also possible.

Definition 2.3.8. Given a vector field $X \in \mathfrak{X}(\mathcal{M})$, streamlines described by Eq. 2.10 are called *integral curves* generated by X.

Proposition 2.3.3. The infinitesimal flow generated by $X \in \mathfrak{X}(\mathcal{M})$ is:

$$x^{\mu}(t) = x^{\mu}(0) + tX^{\mu}(x(t)) + o(t)$$
(2.11)

Definition 2.3.9. A vector field which generates a flow defined for all $t \in \mathbb{R}$ is called *complete*.

Theorem 2.3.3. If \mathcal{M} is compact, then all $X \in \mathfrak{X}(\mathcal{M})$ are complete.

Example 2.3.1. On \mathbb{S}^2 , the flow generated by $X = \partial_{\phi}$ is described by $\dot{\phi} = 1, \dot{\theta} = 0$, thus $\theta(t) = \theta_0$ and $\phi(t) = \phi_0 + t$: the flow lines are lines of constant latitude.

2.3.3 Lie derivative

Defining calculus for vector fields requires a way to compare vectors of different tangent spaces.

Definition 2.3.10. Given a diffeomorphism between two manifolds $\varphi : \mathcal{M} \to \mathcal{N}$ and a function $f : \mathcal{N} \to \mathbb{R}$, the *pull-back* of f is the function $\varphi^* f : \mathcal{M} \to \mathbb{R}$ such that $\varphi^* f(p) = f(\varphi(p))$.

Definition 2.3.11. Given a diffeomorphism between two manifolds $\varphi : \mathcal{M} \to \mathcal{N}$ and a vector field $X \in \mathfrak{X}(\mathcal{M})$, the *push-forward* of X is the vector field $\varphi_*X \in \mathfrak{X}(\mathcal{N})$ such that $\varphi_*X(f) = X(\varphi^*f)$.

This last equality must be evaluated at the appropriate points: $[\varphi_*X(f)](\varphi(p)) = [X(\varphi^*f)](p)$. With the appropriate charts on \mathcal{M} and \mathcal{N} , the definitions above can be rewritten with coordinates:

$$\varphi^* f(x) = f(y(x)) \tag{2.12}$$

$$\varphi_* X(f) = X^{\mu} \frac{\partial f(y(x))}{\partial x^{\mu}} = X^{\mu} \frac{\partial y^{\alpha}}{\partial x^{\mu}} \frac{\partial f(y)}{\partial y^{\alpha}}$$
(2.13)

The notions of pull-back and push-forward allow to compare tangent vectors at neighbouring points and, in particular, to define the derivative along a vector field.

Definition 2.3.12. Given a function $f: \mathcal{M} \to \mathbb{R}$ and a vector field $X \in \mathfrak{X}(\mathcal{M})$, the derivative of f along X (called $Lie\ derivative$) is defined as:

$$\mathcal{L}_X f(x) := \lim_{t \to 0} \frac{f(\sigma_t(x)) - f(x)}{t} = \frac{df(\sigma_t(x))}{dt} \bigg|_{t=0}$$
(2.14)

where σ_t is the flow generated by X.

Proposition 2.3.4. $\mathcal{L}_X f = X(f)$.

Proof.
$$\mathcal{L}_X f = \frac{df(\sigma_t)}{dt} = \frac{\partial f}{\partial x^{\mu}} \frac{dx^{\mu}(t)}{dt} = X^{\mu} \frac{\partial f}{\partial x^{\mu}} = X(f).$$

Definition 2.3.13. Given two vector fields $X, Y \in \mathfrak{X}(\mathcal{M})$, the *Lie derivative* of Y along X is defined as:

$$\mathcal{L}_X Y_p := \lim_{t \to 0} \frac{((\sigma_{-t})_* Y)_p - Y_p}{t}$$
 (2.15)

where σ_t is the flow generated by X.

The use of the inverse flow σ_{-t} is necessary because to evaluate the vector field $\mathcal{L}_X Y$ at the point $p \in \mathcal{M}$, the tangent vector $Y_{\sigma_t(p)} \in T_{\sigma_t(p)} \mathcal{M}$ must be "pushed-back" to $T_p \mathcal{M} = T_{\sigma_0(p)} \mathcal{M}$. With $t \to 0$, the infinitesimal flow σ_{-t} is, according to Eq. 2.11, $x^{\mu}(t) = x^{\mu}(0) - tX^{\mu} + o(t)$, therefore the Lie derivative of base tangent vectors can be expressed as:

$$(\sigma_{-t})_* \partial_\mu = \frac{\partial x^\nu(t)}{\partial x^\mu} \frac{\partial}{\partial x^\nu(t)} = \left(\delta^\nu_\mu - t \frac{\partial X^\nu}{\partial x^\mu} + o(t)\right) \partial_\nu(t) \quad \Longrightarrow \quad \mathcal{L}_X \partial_\mu = -\frac{\partial X^\nu}{\partial x^\mu} \partial_\nu \tag{2.16}$$

Proposition 2.3.5. $\mathcal{L}_X Y = [X, Y]$.

Proof.
$$\mathcal{L}_X Y = \mathcal{L}_X (Y^{\mu} \partial_{\mu}) = (\mathcal{L}_X Y^{\mu}) \partial_{\mu} + Y^{\mu} (\mathcal{L}_X \partial_{\mu}) = X^{\nu} \frac{\partial Y^{\mu}}{\partial x^{\nu}} \partial_{\mu} - Y^{\mu} \frac{\partial X^{\nu}}{\partial x^{\mu}} \partial_{\nu} = [X, Y].$$

Proposition 2.3.6. $\mathcal{L}_X \mathcal{L}_Y Z - \mathcal{L}_Y \mathcal{L}_X Z = \mathcal{L}_{[X,Y]} Z$.

Proof. Trivial with Jacobi identity.

2.4 Tensors

2.4.1 Dual Spaces

Definition 2.4.1. Given a vector space V, its dual V^* is the space of all linear maps $f: V \to \mathbb{R}$.

Given a basis $\{\mathbf{e}_{\mu}\}_{\mu=1,\dots,n}$ of V, its dual basis $\{\mathbf{f}^{\mu}\}_{\mu=1,\dots,n}$ of V^* can be defined by:

$$\mathbf{f}^{\nu}(\mathbf{e}_{\mu}) = \delta^{\nu}_{\mu} \tag{2.17}$$

A general vector in V can be written as $X = X^{\mu} \mathbf{e}_{\mu}$, thus according to Eq. 2.17 $X^{\mu} = \mathbf{f}^{\mu}(X)$.

Proposition 2.4.1. The map $f: \mathbf{e}_{\mu} \mapsto \mathbf{f}^{\mu}$ is an isomorphism between V and V^* .

This isomorphism, however, is basis-dependent.

Proposition 2.4.2. $\dim_{\mathbb{R}} V = \dim_{\mathbb{R}} V^*$.

Proposition 2.4.3. $(V^*)^* = V$.

Proof. The natural isomorphism between $(V^*)^*$ and V is basis-independent: suppose $X \in V$ and $\omega \in V^*$, so that $\omega(X) \in \mathbb{R}$; X can be viewed as $X \in (V^*)^*$ by setting $V(\omega) \equiv \omega(V)$.

2.4.2 Cotangent vectors

Definition 2.4.2. Given a differentiable manifold $(\mathcal{M}, \mathcal{A})$ and a point $p \in \mathcal{M}$, the *cotangent space* to \mathcal{M} at p is defined as $T_p^*\mathcal{M} := (T_p\mathcal{M})^*$.

Elements of $T_p^*\mathcal{M}$ are called *cotangent vectors* (or *covectors*).

Definition 2.4.3. A covector field (or 1-form) is a smooth map $\omega : p \in \mathcal{M} \mapsto \omega_p \in T_p^*\mathcal{M}$. It can also be viewed as a smooth map $\omega : \mathfrak{X}(\mathcal{M}) \to \mathcal{C}^{\infty}(\mathcal{M})$, as $(\omega(X))(p) = \omega_p(X_p) \in \mathbb{R}$.

Definition 2.4.4. The space of all 1-forms on \mathcal{M} is denoted by $\Lambda^1(\mathcal{M})$.

Proposition 2.4.4. $\{dx^{\mu}\}_{\mu=1,\dots,n}$ is a basis of $\Lambda^{1}(\mathcal{M})$ dual to the basis $\{\partial_{\mu}\}_{\mu=1,\dots,n}$ of $\mathfrak{X}(\mathcal{M})$.

Proof. Consider $f \in \mathcal{C}^{\infty}(\mathcal{M})$ and define $df \in \Lambda^{1}(\mathcal{M})$ by df(X) = X(f): taking $f = x^{\mu}$ and $X = \partial_{\mu}$, $df(X) = \partial_{\nu}(x^{\mu}) = \delta^{\mu}_{\nu}$, therefore $\{dx^{\mu}\}_{\mu=1,\dots,n}$ is the dual basis of $\Lambda^{1}(\mathcal{M})$.

This is also confirmed by $df = \frac{\partial f}{\partial x^{\mu}} dx^{\mu}$. These are coordinate basis: in fact, given two different charts $(\varphi, U), (\tilde{\varphi}, V)$:

$$dy^{\mu} = \frac{dy^{\mu}}{dx^{\nu}} dx^{\nu} \tag{2.18}$$

which is the inverse of Eq. 2.3 (not evaluated at a specific point). This ensures that:

$$dy^{\mu} \left(\frac{\partial}{\partial y^{\nu}} \right) = \frac{\partial y^{\mu}}{\partial x^{\alpha}} \frac{\partial x^{\beta}}{\partial y^{\nu}} dx^{\alpha} \left(\frac{\partial}{\partial x^{\beta}} \right) = \frac{\partial y^{\mu}}{\partial x^{\alpha}} \frac{\partial x^{\alpha}}{\partial y^{\nu}} = \delta^{\mu}_{\nu}$$

A 1-form $\omega \in \Lambda^1(\mathcal{M})$ can thus be expressed both as $\omega = \omega_\mu dx^\mu = \tilde{\omega}_\mu dx^\mu$, with:

$$\tilde{\omega}_{\omega} = \frac{\partial x^{\nu}}{\partial y^{\mu}} \omega_{\nu} \tag{2.19}$$

Components of 1-forms which transform this way are called *covariant*.

Definition 2.4.5. Given a diffeomorphism between two manifolds $\varphi : \mathcal{M} \to \mathcal{N}$ and a 1-form $\omega \in \Lambda^1(\mathcal{N})$, the *pull-back* of ω is the 1-form $\varphi^*\omega \in \Lambda^1(\mathcal{M})$ such that $\varphi^*\omega(X) = \omega(\varphi_*X)$.

With the appropriate charts on \mathcal{M} and \mathcal{N} , the definition above can be rewritten with coordinates:

$$\varphi^* \omega = \omega_\alpha \frac{\partial y^\alpha}{\partial x^\mu} dx^\mu \tag{2.20}$$

Definition 2.4.6. Given a vector field $X \in \mathfrak{X}(\mathcal{M})$ and a 1-form $\omega \in \Lambda^1(\mathcal{M})$, the *Lie derivative* of ω along X is defined as:

$$\mathcal{L}_X \omega_p := \lim_{t \to 0} \frac{(\sigma_t^* \omega)_p - \omega_p}{t} \tag{2.21}$$

where σ_t is the flow generated by X.

In contranst with the Lie derivative of a vector field, which pushes forward with σ_{-t} (i.e. pushes back), the Lie derivative of a 1-form pulls back with σ_t : this results in the difference of a minus sign with respect to Eq. 2.16, giving:

$$\mathcal{L}_X dx^\mu = \frac{\partial X^\mu}{\partial x^\nu} dx^\nu \tag{2.22}$$

Therefore, on a general 1-form $\omega = \omega_{\mu} dx^{\mu}$:

$$\mathcal{L}_X \omega = (X^{\nu} \partial_{\nu} \omega_{\mu} + \omega_{\nu} \partial_{\mu} X^{\nu}) dx^{\mu}$$
(2.23)

2.4.3 Tensor fields

Definition 2.4.7. A tensor of rank (r, s) at a $p \in \mathcal{M}$ of a differentiable manifold $(\mathcal{M}, \mathcal{A})$ is a multi-linear map defined as:

$$T_p: \overbrace{T_p^* \mathcal{M} \times \cdots \times T_p^* \mathcal{M}}^r \times \overbrace{T_p \mathcal{M} \times \cdots \times T_p \mathcal{M}}^s \to \mathbb{R}$$
 (2.24)

Example 2.4.1. A cotangent vector $\omega_p \in T_p^* \mathcal{M}$ is a tensor of rank (1,0), while a tangent vector $X_p \in T_p \mathcal{M}$ is a tensor of rank (0,1).

Definition 2.4.8. A tensor field of rank (r, s) is a smooth map $T : p \in \mathcal{M} \mapsto T_p$ tensor of rank (r, s) at p. It can also be viewed as a smooth map $T : [\Lambda^1(\mathcal{M})]^r \times [\mathfrak{X}(\mathcal{M})]^s \to \mathcal{C}^{\infty}(\mathcal{M})$.

Given appropriate basis for vector fields $\{\mathbf{e}_{\mu}\}_{\mu=1,\dots,n}$ and 1-forms $\{\mathbf{f}^{\mu}\}_{\mu=1,\dots,n}$, the components of a tensor field are defined as:

$$T^{\mu_1\dots\mu_r}_{\nu_1\dots\nu_s} := T(\mathbf{f}^{\mu_1},\dots,\mathbf{f}^{\mu_r},\mathbf{e}_{\nu_1},\dots,\mathbf{e}_{\nu_s})$$
 (2.25)

Proposition 2.4.5. On an *n*-dimensional manifold, a (r, s) tensor field has n^{r+s} components, each being element of $\mathcal{C}^{\infty}(\mathcal{M})$.

Consider two general basis transformations, for vector fields and 1-forms, described by invertible matrices A,B such that $\tilde{\mathbf{e}}_{\mu}=A^{\nu}_{\ \mu}\mathbf{e}_{\nu}$ and $\tilde{\mathbf{f}}^{\mu}=B^{\mu}_{\ \nu}\mathbf{f}^{\nu}$, with necessary condition $A^{\mu}_{\ \nu}B^{\rho}_{\ \mu}=\delta^{\rho}_{\ \nu}$ to ensure duality: this implies $B=A^{-1}$, i.e. covectors transform inversely with respect to vectors. Thus:

$$\tilde{T}^{\mu_1 \dots \mu_r}_{\nu_1 \dots \nu_s} = B^{\mu_1}_{\rho_1} \dots B^{\mu_r}_{\rho_r} A^{k_1}_{\nu_1} \dots A^{k_s}_{\nu_s} T^{\rho_1 \dots \rho_r}_{k_1 \dots k_s}$$
(2.26)

If the considered basis are coordinate basis, then $A^{\mu}_{\ \nu} = \frac{\partial x^{\mu}}{\partial y^{\nu}}$ and $B^{\mu}_{\ \nu} = \frac{\partial y^{\mu}}{\partial x^{\nu}}$.

2.4.4 Operations on tensors

Algebric addition and multiplication by functions are trivially defined on tensors of the same rank.

Proposition 2.4.6. The space of all (r, s) tensors at a point $p \in \mathcal{M}$ is denoted by $T_p^{(r,s)}\mathcal{M}$, and it is a vector space.

Definition 2.4.9. Given two tensor fields S of rank (p,q) and T of rank (r,s), their tensor product is defined as:

$$S \otimes T(\omega_1, \dots, \omega_p, \eta_1, \dots, \eta_r, X_1, \dots, X_q, Y_1, \dots, Y_s)$$

$$= S(\omega_1, \dots, \omega_p, X_1, \dots, X_q) T(\eta_1, \dots, \eta_r, Y_1, \dots, Y_s)$$
(2.27)

or, in components:

$$(S \otimes T)^{\mu_1 \dots \mu_p \nu_1 \dots \nu_r}_{\rho_1 \dots \rho_q \dots \sigma_1 \dots \sigma_s} = S^{\mu_1 \dots \mu_p}_{\rho_1 \dots \rho_q} T^{\nu_1 \dots \nu_r}_{\sigma_1 \dots \sigma_s}$$
(2.28)

It is also possible to contract tensor $((r,s) \mapsto (r-1,s-1))$: for example, given a rank (2,1) tensor, a rank (1,0) tensor can be defined as $S(\omega) = T(\omega, \mathbf{f}^{\mu}, \mathbf{e}_{\mu})$, with components $S^{\mu} = T^{\nu\mu}_{\ \mu}$; it must be noted that, in general, $T^{\nu\mu}_{\ \mu} \neq T^{\mu\nu}_{\ \mu}$.

Definition 2.4.10. Given an object $T_{\mu_1...\mu_n}$ dependent on some indices, its *symmetric* and *antisymmetric* parts are respectively defined as:

$$T_{(\mu_1...\mu_n)} := \frac{1}{n!} \sum_{\sigma \in S^n} T_{\sigma(\mu_1)...\sigma(\mu_n)}$$
 (2.29)

$$T_{[\mu_1\dots\mu_n]} := \frac{1}{n!} \sum_{\sigma \in S^n} \operatorname{sgn}(\sigma) T_{\sigma(\mu_1)\dots\sigma(\mu_n)}$$
(2.30)

Conventionally, indices surrounded by || are not (anti-)symmetrized (ex: $T_{[\mu|\nu|\rho]} = \frac{1}{2} (T_{\mu\nu\rho} - T_{\rho\nu\mu})$). As previously seen, vector fields are pushed forward and 1-form are pulled back: tensors will thus behave in a mixed way.

Definition 2.4.11. Given a diffeomorphism between two manifolds $\varphi : \mathcal{M} \to \mathcal{N}$ and a (r, s) tensor field T on \mathcal{M} , the *push-forward* of T is the (r, s) tensor field φ_*T on \mathcal{N} such that, for $\omega_j \in \Lambda^1(\mathcal{N})$ and $X_j \in \mathfrak{X}(\mathcal{N})$:

$$\varphi_* T(\omega_1, \dots, \omega_r, X_1, \dots, X_s) = T(\varphi^* \omega_1, \dots, \varphi^* \omega_r, \varphi_*^{-1} X_1, \dots, \varphi_*^{-1} X_s)$$
(2.31)

Definition 2.4.12. Given a vector field $X \in \mathfrak{X}(\mathcal{M})$ and a (r, s) tensor field T on \mathcal{M} , the *Lie derivative* of T along X is defined as:

$$\mathcal{L}_X T_p := \lim_{t \to 0} \frac{((\sigma_{-t})_* T)_p - T_p}{t}$$
 (2.32)

where σ_t is the flow generated by X.

2.5 Differential forms

Definition 2.5.1. A totally anti-symmetric (0, p) tensor is defined as a *p-form*. The set of all *p*-forms over a manifold \mathcal{M} is denoted as $\Lambda^p(\mathcal{M})$.

Proposition 2.5.1. A p-form has $\binom{n}{p}$ independent components.

Proposition 2.5.2. The maximum degree of differential forms is $p = n \equiv \dim_{\mathbb{R}} \mathcal{M}$: forms in $\Lambda^n(\mathcal{M})$ are called top forms.

Definition 2.5.2. Given $\omega \in \Lambda^p(\mathcal{M}), \eta \in \Lambda^q(\mathcal{M})$, their wedge product is a (p+q)-form defined as:

$$(\omega \wedge \eta)_{\mu_1 \dots \mu_p \nu_1 \dots \nu_q} = \frac{(p+q)!}{p! q!} \omega_{[\mu_1 \dots \mu_p} \eta_{\nu_1 \dots \nu_q]}$$
 (2.33)

Example 2.5.1. Given $\omega, \eta \in \Lambda^2(\mathcal{M})$, their wedge product is $(\omega \wedge \eta)_{\mu\nu} = \omega_{\mu}\eta_{\nu} - \omega_{\nu}\eta_{\mu}$.

Proposition 2.5.3. Given $\omega \in \Lambda^p(\mathcal{M}), \eta \in \Lambda^q(\mathcal{M})$:

$$\omega \wedge \eta = (-1)^{pq} \eta \wedge \omega \tag{2.34}$$

Corollary 2.5.3.1. $\omega \wedge \omega = 0 \ \forall \omega \in \Lambda^p(\mathcal{M}) : p \text{ is odd.}$

Proposition 2.5.4. The wedge product is associative.

Proposition 2.5.5. If $\{\mathbf{f}^{\mu}\}_{\mu=1,\dots,n}$ is a basis of $\Lambda^1(\mathcal{M})$, then $\{\mathbf{f}^{\mu_1} \wedge \dots \wedge \mathbf{f}^{\mu_p}\}_{\mu_1,\dots,\mu_p=1,\dots,n}$ is a basis of $\Lambda^p(\mathcal{M})$.

Locally $\{dx^{\mu}\}_{\mu=1,\dots,n}$ is a basis of $T_p^*\mathcal{M}$, thus a general p-form can be locally written as:

$$\omega = \frac{1}{p!} \omega_{\mu_1 \dots \mu_p} dx^{\mu_1} \wedge \dots \wedge dx^{\mu_p}$$
 (2.35)

Definition 2.5.3. The exterior derivative is a map $d: \Lambda^p(\mathcal{M}) \to \Lambda^{p+1}(\mathcal{M})$ defined as:

$$(d\omega)_{\mu_1...\mu_{p+1}} = (p+1)\partial_{[\mu_1}\omega_{\mu_2...\mu_{p+1}]}$$
(2.36)

In local coordinates:

$$d\omega = \frac{1}{p!} \frac{\partial \omega_{\mu_1 \dots \mu_p}}{\partial x^{\nu}} dx^{\nu} \wedge dx^{\mu_1} \wedge \dots \wedge dx^{\mu_p}$$
(2.37)

Theorem 2.5.1 (Poincaré). $d^2 = 0$.

Proof. Consequence of Schwarz lemma.

Proposition 2.5.6. Given $\omega \in \Lambda^p(\mathcal{M}), \eta \in \Lambda^q(\mathcal{M}),$ then:

$$d(\omega \wedge \eta) = d\omega \wedge \eta + (-1)^p \omega \wedge d\eta \tag{2.38}$$

Proposition 2.5.7. Given a diffeomorphism between two manifolds $\varphi : \mathcal{M} \to \mathcal{N}$ and $\omega \in \Lambda^p(\mathcal{M})$, then $d(\varphi^*\omega) = \varphi^*(d\omega)$.

Corollary 2.5.7.1. Given $X \in \mathfrak{X}(\mathcal{M}), \omega \in \Lambda^p(\mathcal{M})$, then $d(\mathcal{L}_X\omega) = \mathcal{L}_X(d\omega)$.

Definition 2.5.4. $\omega \in \Lambda^p(\mathcal{M})$ is closed if $d\omega = 0$.

Definition 2.5.5. $\omega \in \Lambda^p(\mathcal{M})$ is exact if $\exists \eta \in \Lambda^{p-1}(\mathcal{M}) : \omega = d\eta$.

Theorem 2.5.2. $\omega \in \Lambda^p(\mathcal{M})$ is exact \Rightarrow it is closed.

Definition 2.5.6. Given a vector field $X \in \mathfrak{X}(\mathcal{M})$, the *interior product* determined by X is a map $\iota_X : \Lambda^p(\mathcal{M}) \to \Lambda^{p-1}(\mathcal{M})$ defined as:

$$\iota_X \omega(Y_1, \dots, Y_{p-1}) := \omega(X, Y_1, \dots, Y_{p-1})$$
 (2.39)

On 0-forms (i.e. scalar functions), it is defined as $\iota_X f \equiv 0$.

Proposition 2.5.8. Given $X, Y \in \mathfrak{X}(\mathcal{M})$, then $\iota_X \iota_Y = -\iota_Y \iota_X$.

Proof. Consequence of the total anti-symmetry of p-forms.

Proposition 2.5.9. Given $X \in \mathfrak{X}(\mathcal{M}), \omega \in \Lambda^p(\mathcal{M}), \eta \in \Lambda^q(\mathcal{M}),$ then:

$$\iota_X(\omega \wedge \eta) = \iota_X \omega \wedge \eta + (-1)^p \omega \wedge \iota_X \eta \tag{2.40}$$

Theorem 2.5.3 (Cartan). Given a vector field $X \in \mathfrak{X}(\mathcal{M})$, then:

$$\mathcal{L}_X = d \circ \iota_X + \iota_X \circ d \tag{2.41}$$

Proof. Consider $\omega \in \Lambda^1(\mathcal{M})$:

$$\iota_X(d\omega) = \iota_X \frac{1}{2} \left(\partial_\mu \omega_\nu - \partial_\nu \omega_\mu \right) dx^\mu \wedge dx^\nu = X^\mu \partial_\mu \omega_\nu dx^\nu - X^\nu \partial_\mu \omega_\nu dx^\mu$$

$$d(\iota_X \omega) = d(\omega_\mu X^\mu) = X^\mu \partial_\nu \omega_\mu dx^\nu + \omega_\mu \partial_\nu X^\mu dx^\nu$$

Thus, adding these expressions and recalling Eq. 2.23:

$$(d\iota_X + \iota_X d)\omega = (X^{\mu}\partial_{\mu}\omega_{\nu} + \omega_{\mu}\partial_{\nu}X^{\mu}) dx^{\nu} = \mathcal{L}_X \omega$$

2.5.1 de Rham cohomology

While exact \Rightarrow closed, the converse is not true, in general: it depends on the topological properties of the manifold.

Lemma 2.5.1 (Poincaré). If \mathcal{M} is simply connected, then $\omega \in \Lambda^p(\mathcal{M})$ closed $\Rightarrow \omega$ exact.

In general, it is always possible to choose a simply connected neighbourhood of a point $p \in \mathcal{M}$, in which every closed form is exact, but that may not always be possible globally. It is convenient to set the notation $d_p \equiv d : \Lambda^p(\mathcal{M}) \to \Lambda^{p+1}(\mathcal{M})$.

Definition 2.5.7. The set of all closed p-forms on \mathcal{M} is denoted by $Z^p(\mathcal{M}) := \ker d_p$.

Definition 2.5.8. The set of all exact p-forms on \mathcal{M} is denoted by $B^p(\mathcal{M}) := \operatorname{ran} d_{n-1}$.

Definition 2.5.9. Two closed *p*-forms $\omega, \omega' \in Z^P(\mathcal{M})$ are said to be *equivalent* if $\omega = \omega' + \eta$ for some $\eta \in B^p(\mathcal{M})$.

Definition 2.5.10. The p^{th} de Rham cohomology group of a manifold \mathcal{M} is defined to be:

$$H^p(\mathcal{M}) := Z^p(\mathcal{M})/B^p(\mathcal{M}) \tag{2.42}$$

Definition 2.5.11. The *Betti numbers* of a manifold \mathcal{M} are defined as:

$$B_p := \dim_{\mathbb{R}} H^p(\mathcal{M}) \tag{2.43}$$

Theorem 2.5.4. Given a differentiable manifold, its Betti numbers are always finite.

 $B_0 = 1$ for any connected manifold: there exist constant functions, which are manifestly closed and not exact, due to the non-existence of "-1-forms". Higher Betti numbers are non-zero only if the manifold has some non-trivial topology.

Definition 2.5.12. The *Euler's character* of a manifold \mathcal{M} is defined as:

$$\chi(\mathcal{M}) := \sum_{p \in \mathbb{N}_0} (-1)^p B_p \tag{2.44}$$

Example 2.5.2. The *n*-sphere \mathbb{S}^n has only $B_0 = B_n = 1$, thus $\chi(\mathbb{S}^n) = 1 + (-1)^n$.

Example 2.5.3. The *n*-torus \mathbb{T}^n has $B_p = \binom{n}{p}$, thus $\chi(\mathbb{T}^n) = 0$.

2.5.2 Integration

Definition 2.5.13. A volume form on an n-dimensional differentiable manifold \mathcal{M} is a nowhere-venishing top form v, i.e. locally $v = v(x)dx^1 \wedge \cdots \wedge dx^n : v(x) \neq 0$. If such a form exists, the manifold is said to be *orientable*.

Definition 2.5.14. Given an orientable manifold \mathcal{M} with volume form v, the orientation is:

- right-handed if v(x) > 0 locally on every neighbourhood of \mathcal{M} ;
- left-handed if v(x) < 0 locally on every neighbourhood of \mathcal{M} ;

To ensure that the handedness of the manifold doesn't change on overlapping charts:

$$v = v(x)\frac{\partial x^1}{\partial y^{\mu_1}}dy^{\mu_1} \wedge \dots \wedge \frac{\partial x^n}{\partial x^{\mu_n}}dx^{\mu_n} = v(x)\det\left(\frac{\partial x^\mu}{\partial y^\nu}\right)dy^1 \wedge \dots \wedge dy^n$$

It is therefore necessary that the two sets of coordinates on the overlapping region satisfy:

$$\det\left(\frac{\partial x^{\mu}}{\partial y^{\nu}}\right) > 0 \tag{2.45}$$

Non-orientable manifolds cannot be covered by overlapping charts satisfying this condition.

Example 2.5.4. The real projective space \mathbb{RP}^n is orientable for odd n and non-orientable for even n. Example 2.5.5. The complex projective space \mathbb{CP}^n is orientable for all $n \in \mathbb{N}$. **Definition 2.5.15.** Given a function $f: \mathcal{M} \to \mathbb{R}$ on an orientable manifold \mathcal{M} with volume form v and a chart (φ, U) on \mathcal{M} with coordinates $\{x^{\mu}\}_{\mu=1,\dots,n}$, the *integral* of f on $O = \varphi^{-1}(U) \subset \mathcal{M}$ is defined as:

 $\int_{O} fv := \int_{U} dx_1 \dots dx_n f(x)v(x) \tag{2.46}$

It is clear that the volume form acts like a measure on the manifold. To integrate over the whole manifold, it must be divided up into different regions, each covered by a single chart.

Definition 2.5.16. A k-dimensional manifold Σ is a submanifold of an n-dimensionale manifold \mathcal{M} , with n > k, if there exists an injective map $\varphi : \Sigma \to \mathcal{M}$ such that $\varphi_* : T_p(\Sigma) \to T_{\varphi(p)}(\mathcal{M})$ is injective.

Definition 2.5.17. Given a k-form $\omega \in \Lambda^k(\mathcal{M})$, its integral over a k-dimensional submanifold Σ of \mathcal{M} is defined as:

$$\int_{\varphi(\Sigma)} \omega := \int_{\Sigma} \varphi^* \omega \tag{2.47}$$

Example 2.5.6. Consider a 1-form $\omega \in \Lambda^1(\mathcal{M})$ and a 1-dimensional submanifold γ of \mathcal{M} described by a curve $\sigma : \gamma \to \mathcal{M} : x^{\mu} = \sigma^{\mu}(t)$: locally $\omega = \omega_{\mu}(x)dx^{\mu}$, thus the integral of ω on γ can be calculated as $\int_{\sigma(\gamma)} \omega = \int_{\gamma} \sigma^* \omega = \int_{\gamma} d\tau \, \omega_{\mu}(x) \frac{dx^{\mu}}{d\tau}$.

2.5.2.1 Stokes' theorem

Integration can be generalized beyond smooth (i.e. differentiable) manifolds.

Definition 2.5.18. An *n*-dimensional manifold with boundary is a Hausdorff topological space, equipped with a compatible maximal atlas, which is locally homeomorphic to $\mathbb{R}^{n-1} \times [a, \infty) : a \in \mathbb{R}$. The boundary $\partial \mathcal{M}$ is the 1-dimensional submanifold determined by $x^n = a$.

Theorem 2.5.5 (Stokes). Given an n-dimensional manifold \mathcal{M} with boundary $\partial \mathcal{M}$, then for any $\omega \in \Lambda^{n-1}(\mathcal{M})$:

$$\int_{\mathcal{M}} d\omega = \int_{\partial \mathcal{M}} \omega \tag{2.48}$$

This important theorem unifies many different results.

Given the 1-dimensional manifold $I = [a, b] \subset \mathbb{R}$, then for any 0-form (i.e. scalar function) $\omega = \omega(x)$:

$$\int_{I} d\omega = \int_{a}^{b} \frac{d\omega}{dx} dx = \int_{\partial I} \omega = \omega(b) - \omega(a)$$

which is the fundamental theorem of calculus.

Given a 2-dimensional manifold with boundary $S \subset \mathbb{R}^2$ and a 1-form $\omega = \omega_1 dx^1 + \omega_2 dx^2$, then $d\omega = (\partial_1 \omega_2 - \partial_2 \omega_1) dx^1 \wedge dx^2$ and:

$$\int_{S} d\omega = \int_{S} \left(\frac{\partial \omega_{2}}{\partial x^{1}} - \frac{\partial \omega_{1}}{\partial x^{2}} \right) dx^{1} dx^{2} = \int_{\partial S} \omega = \int_{\partial S} \omega_{1} dx^{1} + \omega_{2} dx^{2}$$

which is Green's theorem.

GIven a 3-dimensional manifold with boundary $V \subset \mathbb{R}^3$ and a 2-form $\omega = \omega_1 dx^2 \wedge dx^3 + \omega_2 dx^3 \wedge dx^1 + \omega_3 dx^1 \wedge dx^2$, then $d\omega = (\partial_1 \omega_1 + \partial_2 \omega_2 + \partial_3 \omega_3) dx^1 \wedge dx^2 \wedge dx^3$ and:

$$\int_{V} d\omega = \int_{V} \left(\frac{\partial \omega_{1}}{\partial x^{1}} + \frac{\partial \omega_{2}}{\partial x^{2}} + \frac{\partial \omega_{3}}{\partial x^{3}} \right) = \int_{\partial V} \omega = \int_{\partial V} \omega_{1} dx^{2} dx^{3} + \omega_{2} dx^{3} dx^{1} + \omega_{3} dx^{1} dx^{2}$$

which is Gauss' theorem.

Riemannian Geometry

3.1 Metric manifolds

Definition 3.1.1. A metric g is a (0,2) tensor field on a manifold \mathcal{M} that is:

- 1. symmetric: g(X,Y) = g(Y,X);
- 2. non-degenerate: $\exists p \in \mathcal{M} : g(X,Y)|_p = 0 \, \forall Y \in T_p \mathcal{M} \Rightarrow X_p = 0.$

Definition 3.1.2. A metric manifold (\mathcal{M}, g) is a manifold equipped with a metric.

With a choice of coordinates, the metric can be written as:

$$g = g_{\mu\nu}(x)dx^{\mu} \otimes dx^{\nu} \tag{3.1}$$

where:

$$g_{\mu\nu} = g\left(\frac{\partial}{\partial x^{\mu}}, \frac{\partial}{\partial x^{\nu}}\right) \tag{3.2}$$

It is often written also as $ds^2 = g_{\mu\nu}(x)dx^{\mu}dx^{\nu}$. The matrix $g_{\mu\nu}(x) \in \mathbb{R}^{n\times n}$ is symmetric, and there's always a choice of basis on each tangent space such that this matrix is diagonal: the non-degeneracy condition implies that none of the diagonal elements vanish.

Proposition 3.1.1. The signature of a metric, i.e. the number of negative entries when diagonalized, is independent on the choice of basis.

Proof. From Sylvester's theorem of inertia.

Riemannian manifolds

Definition 3.1.3. A Riemannian manifold (\mathcal{M}, g) is a manifold equipped with a metric with totally-positive signature.

Example 3.1.1. The Euclidean space \mathbb{R}^n , equipped with the metric $g_{\mu\nu} = \delta_{\mu\nu}$ (in Cartesian coordinates), is a Riemannian manifold.

Definition 3.1.4. Given a Riemannian manifold (\mathcal{M}, g) and $X \in \mathfrak{X}(\mathcal{M})$, the *length* of X at $p \in \mathcal{M}$ is:

$$|X_p| := \sqrt{g(X,X)|_p} \tag{3.3}$$

Given $Y \in \mathfrak{X}(\mathcal{M})$, the angle between X and Y at $p \in \mathcal{M}$ is:

$$\cos \theta := \frac{g(X,Y)|_p}{|X_p| |Y_p|}$$
 (3.4)

This can be generalized to distances between points on a curve $\sigma: \mathbb{R} \to \mathcal{M}$:

$$d(p,q) = \int_a^b dt \sqrt{g(X,X)|_{\sigma(t)}}$$
(3.5)

where $\sigma(a) = p$, $\sigma(b) = q$ and X is the tangent vector field of the curve. With parametrization $x^{\mu}(t)$, the tangent vector has components $X^{\mu} = \frac{dx^{\mu}}{dt}$, thus:

$$d(p,q) = \int_{a}^{b} dt \sqrt{g_{\mu\nu}(x) \frac{dx^{\mu}}{dt} \frac{dx^{\nu}}{dt}}$$
(3.6)

It is important to note that this distance is independent of the parametrization.

Lorentzian manifolds

Definition 3.1.5. A Lorentzian manifold (\mathcal{M}, g) is a manifold equipped with a metric which has a signature with a single negative sign.

Example 3.1.2. The simplest Lorentzian manifold is \mathbb{R}^n with the Minkowski metric:

$$\eta = -dx^0 \otimes dx^0 + dx^1 \otimes dx^1 + \dots + dx^{n-1} \otimes dx^{n-1}$$
(3.7)

Its components are $\eta_{\mu\nu} = \text{diag}(-1, +1, \dots, +1)$, thus this is a Lorentzian manifold.

On a general Lorentzian manifold, at any point $p \in \mathcal{M}$ it is always possible to choose an orthonormal basis $\{e_{\mu}\}_{\mu=0,\dots,n-1}$ of $T_{p}\mathcal{M}$ such that $g_{\mu\nu}|_{p}=\eta_{\mu\nu}$: this fact is closely related to the equivalence principle. Consider a different basis $\tilde{e}_{\mu}=\Lambda^{\nu}_{\mu}e_{\nu}$: the condition for it to leave the Minkowski metric unchanged is:

$$\eta_{\mu\nu} = \Lambda^{\rho}_{\ \mu} \Lambda^{\sigma}_{\ \nu} \eta_{\rho\sigma} \tag{3.8}$$

This is the defining equation of a Lorentz transformation: on a Lorentzian manifold, the basic features of special relativity are locally recovered. Thus, other ideas from special relativity can be imported.

Definition 3.1.6. Given a Lorentzian manifold (\mathcal{M}, g) and $X \in \mathfrak{X}(\mathcal{M})$, at $p \in \mathcal{M}$ the vector field is said to be:

- $timelike if g(X_p, X_p) < 0;$
- $null \text{ if } g(X_p, X_p) = 0;$
- spacelike if $q(X_n, X_n) > 0$.

At each point $p \in \mathcal{M}$ it is possible to draw *lightcones*, i.e. the null tangent vectors at that point, which are past-directed or future-directed: these lightcones vary smoothly as the point is varied smoothly on the manifold, elucidating the causal structure of spacetime.

The distance between two points on a curve depends on the nature of the tangent vector field of the curve: a *timelike curve* is a curve whose tangent vector field is everywhere timelike, and analogously for the other cases. The distance on a spacelike curve is defined as in Eq. 3.5, while that on a timelike curve gets a negative sign in the square root. With parametrization $x^{\mu}(t)$, it is possible to define the *proper time* on a timelike curve as:

$$\tau = \int_{a}^{b} dt \sqrt{-g_{\mu\nu}(x) \frac{dx^{\mu}}{dt} \frac{dx^{\nu}}{dt}}$$
 (3.9)

This is precisely the action of a free particle moving in spacetime.

3.1.1 Metric properties

The metric defines a natural isomorphism between vectors and covectors.

Proposition 3.1.2. Given a metric manifold (\mathcal{M}, g) , the metric defines for each $p \in \mathcal{M}$ a natural isomorphism $g: X_p \in T_p \mathcal{M} \to \omega_p \in T_p^* \mathcal{M}: \omega_p(Y_p) = g(X_p, Y_p) \, \forall Y_p \in T_p \mathcal{M}$.

In a chosen coordinate basis, the vector $X = X^{\mu}\partial_{\mu}$ is mapped to the one-form $X = X_{\mu}dx^{\mu}$, thus the following identity holds:

$$X_{\mu} = g_{\mu\nu} X^{\nu} \tag{3.10}$$

Being g non-degenerate, the matrix $g_{\mu\nu}$ is invertible, with inverse $g^{\mu\nu}$ such that:

$$g^{\mu\nu}g_{\nu\rho} = \delta^{\mu}_{\rho} \tag{3.11}$$

Its elements are the components of a (2,0) symmetric tensor $\hat{g} := g^{\mu\nu} \partial_{\mu} \otimes \partial_{\nu}$, which defines the inverse of the natural isomorphism in Prop. 3.1.2:

$$X^{\mu} = g^{\mu\nu} X_{\nu} \tag{3.12}$$

The metric also defines a natural volume form on the manifold.

Definition 3.1.7. Given an *n*-dimensional metric manifold (\mathcal{M}, g) , the *volume form* is the top-form:

$$v := \sqrt{g} \, dx^1 \wedge \dots \wedge dx^n \tag{3.13}$$

where $g := |\det g_{\mu\nu}|$.

Proposition 3.1.3. The volume form is basis-independent.

Proof. Consider a new set of coordinates y^{μ} such that $dx^{\mu} = A^{\mu}_{\ \nu} dy^{\nu}$, where $A^{\mu}_{\ \nu} = \frac{\partial x^{\mu}}{\partial y^{\nu}}$. In general:

$$dx^1 \wedge \dots \wedge dx^n = A^1_{\mu_1} \dots A^n_{\mu_n} dy^{\mu_1} \wedge \dots \wedge dy^{\mu_n}$$

Recalling the anti-symmetry of the wedge product and the definition of determinant, this can be rewritten as:

$$dx^{1} \wedge \dots \wedge dx^{n} = \sum_{\pi \in S^{n}} \operatorname{sgn} \pi A^{1}_{\pi(1)} \dots A^{n}_{\pi(n)} dy^{1} \wedge \dots \wedge dy^{n} = \det A \, dy^{1} \wedge \dots \wedge dy^{n}$$

Note the Jacobian factor which arises when changing the measure. On the other hand:

$$g_{\mu\nu} = \frac{\partial y^{\rho}}{\partial x^{\mu}} \frac{\partial y^{\sigma}}{\partial x^{\nu}} \tilde{g}_{\rho\sigma} = (A^{-1})^{\rho}_{\ \mu} (A^{-1})^{\sigma}_{\ \nu} \tilde{g}_{\rho\sigma} \quad \Rightarrow \quad \det g_{\mu\nu} = \frac{\det \tilde{g}_{\mu\nu}}{(\det A)^2}$$

The factors $\det A$ and $(\det A)^{-1}$ cancel, thus yielding the thesis.

The volume form can be rewritten as:

$$v = \frac{1}{n!} v_{\mu_1 \dots \mu_n} dx^{\mu_1} \wedge \dots \wedge dx^{\mu_n} \equiv \frac{1}{n!} \sqrt{g} \,\epsilon_{\mu_1 \dots \mu_n} dx^{\mu_1} \wedge \dots \wedge dx^{\mu_n}$$
(3.14)

where $\epsilon_{\mu_1...\mu_n}$ is the totally-antisymmetric *n*-dimensional symbol (generalization of the Levi-Civita symbol). $\epsilon_{\mu_1...\mu_n}$ cannot be considered a proper tensor, as its components are always +1, -1, 0

indipendently if the indices are covariant or contravariant: it is, in fact, a tensor density, i.e. a tensor divided by \sqrt{g} . It can be shown that:

$$v^{\mu_1...\mu_n} = g^{\mu_1\nu_1} \dots g^{\mu_n\nu_n} v_{\mu_1...\mu_n} = \sigma \frac{1}{\sqrt{g}} \epsilon^{\mu_1...\mu_n}$$
(3.15)

where σ is the sign of the signature (ex.: $\sigma = +1$ for Riemannian manifolds and $\sigma = -1$ for Lorentzian manifolds). As notation, the integral of a generic function f on \mathcal{M} is denoted as:

$$\int_{\mathcal{M}} f v \equiv \int_{\mathcal{M}} d^n x \sqrt{g} f \tag{3.16}$$

3.1.1.1 Hodge theory

Definition 3.1.8. Given an *n*-dimensional oriented metric manifold (\mathcal{M}, g) , the *Hodge dual* is defined as the map $\star : \Lambda^p(\mathcal{M}) \to \Lambda^{n-p}(\mathcal{M}) : \omega \mapsto \star \omega$ such that:

$$\star \omega_{\mu_1 \dots \mu_{n-p}} := \frac{1}{(n-p)!} \sqrt{g} \,\epsilon_{\mu_1 \dots \mu_{n-p} \nu_1 \dots \nu_p} \omega^{\nu_1 \dots \nu_p} \tag{3.17}$$

In this section, the orientedness and n-dimensionality of the manifold are implied.

Proposition 3.1.4. The Hodge dual is basis-independent.

It is useful to state a lemma for future calculations.

Lemma 3.1.1.
$$v^{\mu_1...\mu_p\rho_1...\rho_{n-p}}v_{\nu_1...\nu_p\rho_1...\rho_{n-p}} = \sigma p!(n-p)!\delta^{\mu_1}_{[\nu_1}...\delta^{\mu_p]}$$
.

Proposition 3.1.5. $\star(\star\omega) = \sigma(-1)^{p(n-p)}\omega$.

The Hodge dual defines an inner product on each $\Lambda^p(\mathcal{M})$:

$$\langle \omega, \eta \rangle := \int_{\mathcal{M}} \omega \wedge \star \eta \tag{3.18}$$

This allows to define operators and their adjoints on the form spaces.

Proposition 3.1.6. Given a metric manifold (\mathcal{M}, g) and two forms $\omega \in \Lambda^p(\mathcal{M}), \alpha \in \Lambda^{p-1}(\mathcal{M})$, then:

$$\langle d\alpha, \omega \rangle = \langle \alpha, d^{\dagger}\omega \rangle \tag{3.19}$$

where the adjoint of the exterior derivative $d^{\dagger}: \Lambda^{p}(\mathcal{M}) \to \Lambda^{p-1}(\mathcal{M})$ is defined as:

$$d^{\dagger} := \sigma \left(-1\right)^{np+n-1} \star d \star \tag{3.20}$$

Proof. To simplify the proof, consider a closed manifold; then, from Stokes' theorem and Eq. 2.38:

$$0 = \int_{\mathcal{M}} d(\alpha \wedge \star \omega) = \langle d\alpha, \omega \rangle + \int_{\mathcal{M}} (-1)^{p-1} \alpha \wedge d \star \omega$$

The second term is proportional to $\langle \alpha, \star d \star \omega \rangle$: to determine the relative sign, note that $d \star \omega \in \Lambda^{n-p+1}(\mathcal{M})$, thus, from Prop. 3.1.5, $\star \star d \star \omega = \sigma (-1)^{(n-p+1)(p-1)} d \star \omega$. In conclusion:

$$\langle \alpha, \star \, d \star \omega \rangle = \sigma \, (-1)^{(n-p)(p-1)} \int_{\mathcal{M}} (-1)^{p-1} \, \alpha \wedge \, d \star \omega \quad \Rightarrow \quad \langle d\alpha, \omega \rangle = \sigma \, (-1)^{(n-p)(p-1)+1} \, \langle \alpha, \star \, d \star \omega \rangle$$

Noting that $(-1)^{(n-p)(p-1)+1} = (-1)^{np+n-1}$, as in general $(-1)^{-n} = (-1)^n$ and $(-1)^{-p^2+p+1} = (-1)^{-1}$ due to p(p-1) being always even, concludes the proof.

Definition 3.1.9. Given a metric manifold (\mathcal{M}, g) , the Laplacian $\triangle : \Lambda^p(\mathcal{M}) \to \Lambda^p(\mathcal{M})$ is defined as the operator:

$$\Delta := (d + d^{\dagger})^2 \tag{3.21}$$

Proposition 3.1.7. $\triangle = dd^{\dagger} + d^{\dagger}d = \{d, d^{\dagger}\}.$

Proof. Trivial, given
$$d^2 = d^{\dagger 2} = 0$$
.

It is possible to calculate an explicit expression for the Laplacian of functions.

Lemma 3.1.2. Given $f \in \mathcal{C}^{\infty}(\mathcal{M})$, then $d^{\dagger}f = 0$.

Proof. Trivial noting that $\star f$ is a top-form.

Proposition 3.1.8. Given $f \in C^{\infty}(\mathcal{M})$, then:

$$\Delta f = -\frac{\sigma}{\sqrt{g}} \partial_{\nu} \left(\sqrt{g} g^{\mu\nu} \partial_{\mu} f \right) \tag{3.22}$$

Proof. Via direct calculation, using Lemma 3.1.2:

$$\Delta f = \sigma (-1)^{n^2+n-1} \star d \star (\partial_{\mu} f dx^{\mu}) = -\sigma \star d (\partial_{\mu} f \star dx^{\mu})
= -\frac{\sigma}{(n-1)!} \star d (\partial_{\mu} f g^{\mu\nu} \sqrt{g} \epsilon_{\nu\rho_1...\rho_{n-1}} dx^{\rho_1} \wedge \cdots \wedge dx^{\rho_{n-1}})
= -\frac{\sigma}{(n-1)!} \star \partial_{\alpha} (\sqrt{g} g^{\mu\nu} \partial_{\mu} f) \epsilon_{\nu\rho_1...\rho_{n-1}} dx^{\alpha} \wedge dx^{\rho_1} \wedge \cdots \wedge dx^{\rho_{n-1}}
= -\sigma \star \partial_{\nu} (\sqrt{g} g^{\mu\nu} \partial_{\mu} f) dx^1 \wedge \ldots dx^n = -\frac{\sigma}{\sqrt{g}} \partial_{\nu} (\sqrt{g} g^{\mu\nu} \partial_{\mu} f)$$

The Laplacian operator is linked to the de Rham cohomology.

Definition 3.1.10. Given $\omega \in \Lambda^p(\mathcal{M})$, it is said to be harmonic if $\Delta \omega = 0$.

Definition 3.1.11. The space of harmonic p-forms on (\mathcal{M}, g) is denoted as $\operatorname{Harm}^p(\mathcal{M})$.

Proposition 3.1.9. A harmonic form is both closed and co-closed.

Proof. $0 = \langle \omega, \triangle \omega \rangle = \langle d\omega, d\omega \rangle + \langle d^{\dagger}\omega, d^{\dagger}\omega \rangle$, thus $d\omega = 0$ and $d^{\dagger}\omega = 0$, for the inner product is positive-defined.

Theorem 3.1.1. Given a compact Riemannian manifold (\mathcal{M}, g) , any $\omega \in \Lambda^p(\mathcal{M})$ can be uniquely decomposed as $\omega = d\alpha + d^{\dagger}\beta + \gamma$, with $\alpha \in \Lambda^{p-1}(\mathcal{M})$, $\beta \in \Lambda^{p+1}(\mathcal{M})$ and $\gamma \in \operatorname{Harm}^p(\mathcal{M})$.

Theorem 3.1.2 (Hodge). Given a compact Riemannian manifold (\mathcal{M}, g) , there is an isomorphism:

$$\operatorname{Harm}^p(\mathcal{M}) \cong H^p(\mathcal{M})$$
 (3.23)

Proof. From Prop. 3.1.9 $\operatorname{Harm}^p(\mathcal{M}) \subset Z^p(\mathcal{M})$, but the uniqueness of decomposition in Th. 3.1.1 implies $\forall \gamma \in \operatorname{Harm}^p(\mathcal{M}) \exists \eta_{\gamma} \in \Lambda^{p-1}(\mathcal{M}) : \gamma \neq d\eta_{\gamma}$, thus $\operatorname{Harm}^p(\mathcal{M}) \subset H^p(\mathcal{M})$.

WTS that any equivalence class $[\omega] \in H^p(\mathcal{M})$ can be represented by a harmonic form. By Th, 3.1.1 $\omega = d\alpha + d^{\dagger}\beta + \gamma$, but $\omega \in H^p(\mathcal{M})$ implies $d\omega = 0$ by definition, so:

$$0 = \langle d\omega, \beta \rangle = \langle \omega, d^{\dagger}\beta \rangle = \langle d\alpha + d^{\dagger}\beta + \gamma, d^{\dagger}\beta \rangle = \langle d^{\dagger}\beta, d^{\dagger}\beta \rangle$$

The inner product is positive-definite, thus $d^{\dagger}\beta = 0$, hence $\omega = \gamma + d\alpha$. By definition $H^p(\mathcal{M}) := Z^p(\mathcal{M})/B^p(\mathcal{M})$, so $[\omega] = \gamma$.

Corollary 3.1.2.1. $B_p = \dim_{\mathbb{R}} \operatorname{Harm}^p(\mathcal{M})$.

3.2 Connections

There's a different way to differentiate tensor fields distinct from the Lie derivative, associated to a different way to map different vector spaces at different points: the covariant derivative. From now on, \mathcal{M} is implied to be an n-dimensional metric manifold with metric q.

3.2.1 Covariant derivative

Definition 3.2.1. The connection is a map $\nabla: \mathfrak{X}(\mathcal{M}) \times \mathfrak{X}(\mathcal{M}) \to \mathfrak{X}(\mathcal{M})$, usually written as $\nabla(X,Y) \equiv \nabla_X Y$, where ∇_X is called the *covariant derivative*, satisfying the following properties for all $X,Y,Z \in \mathfrak{X}(\mathcal{M})$:

- 1. $\nabla_X(Y+Z) = \nabla_X Y + \nabla_X Z;$
- 2. $\nabla_{fX+gY}Z = f\nabla_XZ + g\nabla_YZ\,\forall f,g\in\mathcal{C}^{\infty}(\mathcal{M});$
- 3. $\nabla_X(fY) = f\nabla_X Y + X(f)Y \,\forall f \in \mathcal{C}^{\infty}(\mathcal{M}).$

Usually $X(f) \equiv \nabla_X f$. The covariant derivative endows the manifold with more structure: in particular, given a basis $\{e_{\mu}\}$ of $\mathfrak{X}(\mathcal{M})$, its covariant derivative is expressed as:

$$\nabla_{e_{\rho}} e_{\mu} \equiv \Gamma^{\mu}_{\rho\nu} e_{\mu} \tag{3.24}$$

The $\Gamma^{\mu}_{\rho\nu}$ are the components of the connection on that basis. Usually $\nabla_{e_{\mu}} \equiv \nabla_{\mu}$, thus resembling a partial derivative. To elucidate how the covariant derivative acts on vector fields:

$$\nabla_{X}Y = \nabla_{X}(Y^{\mu}e_{\mu})
= X(Y^{\mu})e_{\mu} + Y^{\mu}\nabla_{X}e_{\mu}
= X^{\nu}e_{\nu}(Y^{\mu})e_{\mu} + Y^{\mu}X^{\nu}\nabla_{\nu}e_{\mu}
= X^{\nu}\left[e_{\nu}(Y^{\mu}) + \Gamma^{\mu}_{\nu\rho}Y^{\rho}\right]e_{\mu}
= X^{\nu}\nabla_{\nu}Y = X^{\nu}(\nabla_{\nu}Y)^{\mu}e_{\mu}$$

The dependency on X can therefore be eliminated, and in components:

$$(\nabla_{\nu}Y)^{\mu} = e_{\nu}(Y^{\mu}) + \Gamma^{\mu}_{\nu\rho}Y^{\rho} \tag{3.25}$$

A sloppy notation is often used: $(\nabla_{\nu}Y)^{\mu} \equiv \nabla_{\nu}Y^{\mu}$. This must not be confused as the covariant derivative of Y^{μ} . Moreover $\nabla_{\nu}Y^{\mu} \equiv Y^{\mu}_{:\nu}$, while $\partial_{\mu}f \equiv f_{,\mu}$. On the coordinate basis $e_{\mu} = \partial_{\mu}$, then:

$$Y^{\mu}_{;\nu} = Y^{\mu}_{,\nu} + \Gamma^{\mu}_{\nu\rho} Y^{\rho} \tag{3.26}$$

Note that $Y^{\mu}_{;\nu}$ is the μ^{th} component of $\nabla_{\nu}Y$, while $Y^{\mu}_{,\nu}$ is the partial derivative of Y^{μ} along ∂_{ν} . The covariant derivative coincides with other derivatives on $\mathcal{C}^{\infty}(\mathcal{M})$: it can be shown that $\nabla_X f = \mathcal{L}_X f = X(f)$ and $\nabla_{\mu} f = \partial_{\mu} f$. On $\mathfrak{X}(\mathcal{M})$, however, ∇_X and \mathcal{L}_X are distinct: while $\nabla_X = X^{\mu} \nabla_{\mu}$, there's no way to write the same relation for \mathcal{L}_X , for it depends not only on X but on its first derivative too. The covariant derivative is thus the natural generalization of the partial derivative to curved manifolds.

Proposition 3.2.1. $\Gamma^{\mu}_{\rho\nu}$ are not components of a tensor.

Proof. Given the basis transformation $\tilde{e}_{\nu} = A^{\mu}_{\nu} e_{\mu}$, with A and invertible matrix (if they're both coordinate basis, then $A^{\mu}_{\nu} = \frac{\partial x^{\mu}}{\partial x^{\nu}}$), the components of a (1,2) tensor must transform as:

$$\tilde{T}^{\mu}_{\ \rho\nu} = (A^{-1})^{\mu}_{\ \tau} A^{\sigma}_{\ \rho} A^{\lambda}_{\ \nu} T^{\tau}_{\ \sigma\lambda}$$

In the new basis:

$$\begin{split} \tilde{\Gamma}^{\mu}_{\rho\nu}\tilde{e}_{\mu} &= \nabla_{\tilde{e}_{\rho}}\tilde{e}_{\nu} = \nabla_{A^{\sigma}_{\rho}e_{\sigma}}(A^{\lambda}_{\nu}e_{\nu}) = A^{\sigma}_{\rho}\nabla_{e_{\sigma}}(A^{\lambda}_{\nu}e_{\lambda}) \\ &= A^{\sigma}_{\rho}A^{\lambda}_{\nu}\Gamma^{\tau}_{\sigma\lambda}e_{\tau} + A^{\sigma}_{\rho}e_{\lambda}\partial_{\sigma}A^{\lambda}_{\nu} = \left[A^{\sigma}_{\rho}A^{\lambda}_{\nu}\Gamma^{\tau}_{\lambda} + A^{\sigma}_{\rho}\partial_{\sigma}A^{\tau}_{\nu}\right]e_{\tau} \\ &= \left[A^{\sigma}_{\rho}A^{\lambda}_{\nu}\Gamma^{\tau}_{\sigma\lambda} + A^{\sigma}_{\rho}\partial_{\sigma}A^{\tau}_{\nu}\right](A^{-1})^{\mu}_{\tau}\tilde{e}_{\mu} \end{split}$$

Thus, there's a second term proportional to ∂A which deviates from the transformation law:

$$\tilde{\Gamma}^{\mu}_{\rho\nu} = (A^{-1})^{\mu}_{\ \tau} A^{\sigma}_{\ \rho} \left[A^{\lambda}_{\ \nu} \Gamma^{\tau}_{\sigma\lambda} + \partial_{\sigma} A^{\tau}_{\ \nu} \right]$$

3.2.2 Covariant derivative of tensors

First of all, it is necessary to elucidate how the covariant derivative acts on one-forms. Given a one-form ω , the one-form $\nabla_X \omega$ is defined by its action on vector fields. By Leibniz rule:

$$\nabla_X(\omega(Y)) = (\nabla_X \omega)(Y) + \omega(\nabla_X Y)$$

Recalling that $\omega(Y)$ is a function, $\nabla_X(\omega(Y)) = X(\omega(Y))$, therefore:

$$(\nabla_X \omega)(Y) = X(\omega(Y)) - \omega(\nabla_X Y) \tag{3.27}$$

Expressing it in coordinates:

$$X^{\mu}(\nabla_{\mu}\omega)_{\nu}Y^{\nu} = X^{\mu}\partial_{\mu}(\omega_{\nu}Y^{\nu}) - \omega_{\nu}X^{\mu} \left[\partial_{\mu}Y^{\nu} + \Gamma^{\nu}_{\mu\rho}Y^{\rho}\right]$$
$$= X^{\mu} \left[\partial_{\mu}\omega_{\rho} - \Gamma^{\nu}_{\mu\rho}\omega_{\nu}\right]Y^{\rho}$$

Crucially, the ∂Y terms cancel out, allowing to define $\nabla_X \omega$ without referencing Y:

$$(\nabla_{\mu}\omega)_{\rho} = \partial_{\mu}\omega_{\rho} - \Gamma^{\nu}_{\mu\rho}\omega_{\nu} \tag{3.28}$$

Using the same notation as for vector fields $(\nabla_{\mu}\omega)_{\rho} \equiv \nabla_{\mu}\omega_{\rho} \equiv \omega_{\rho;\mu}$:

$$\omega_{\rho;\mu} = \omega_{\rho,\mu} - \Gamma^{\nu}_{\mu\rho}\omega_{\nu} \tag{3.29}$$

This kind of argument can be extended to a general (p,q) tensor field:

$$\nabla_{\rho} T^{\mu_{1}\dots\mu_{p}}_{\nu_{1}\dots\mu_{q}} = \partial_{\rho} T^{\mu_{1}\dots\mu_{p}}_{\nu_{1}\dots\nu_{q}} + \Gamma^{\mu_{1}}_{\rho\sigma} T^{\sigma\mu_{2}\dots\mu_{p}}_{\nu_{1}\dots\nu_{q}} + \dots + \Gamma^{\mu_{p}}_{\rho\sigma} T^{\mu_{1}\dots\mu_{p-1}\sigma}_{\nu_{1}\dots\nu_{q}} - \Gamma^{\sigma}_{\rho\nu_{1}} T^{\mu_{1}\dots\mu_{p}}_{\sigma\nu_{2}\dots\nu_{q}} - \dots - \Gamma^{\sigma}_{\rho\nu_{q}} T^{\mu_{1}\dots\mu_{p}}_{\nu_{1}\dots\nu_{q-1}\sigma}$$
(3.30)

The pattern is clear: for each upper index μ there's a $+\Gamma^{\mu}_{\rho\sigma}T^{\sigma}$ term, while for each lower index ν there's $-\Gamma^{\sigma}_{\rho\nu}T_{\sigma}$ term. Furthermore, it is necessary to generalize the comma-notation: for example, $X^{\mu}_{;\nu\rho} \equiv \nabla_{\rho}\nabla_{\nu}X^{\mu}$, so the rightmost index is the one whose covariant derivative acts first.

3.2.2.1 Torsion and curvature

Even though the connection is not a tensor, it is used to construct two important tensors.

Definition 3.2.2. The torsion is a (1,2) tensor defined on $\Lambda^1(\mathcal{M}) \times \mathfrak{X}(\mathcal{M}) \times \mathfrak{X}(\mathcal{M})$ as:

$$T(\omega, X, Y) := \omega(\nabla_X Y - \nabla_Y X - [X, Y]) \tag{3.31}$$

Alternatively, the torsion can be viewed as a map $T: \mathfrak{X}(\mathcal{M}) \times \mathfrak{X}(\mathcal{M}) \to \mathfrak{X}(\mathcal{M})$ such that:

$$T(X,Y) = \nabla_X Y - \nabla_Y X - [X,Y] \tag{3.32}$$

Definition 3.2.3. The *curvature* is a (1,3) tensor defined on $\Lambda^1(\mathcal{M}) \times \mathfrak{X}(\mathcal{M}) \times \mathfrak{X}(\mathcal{M}) \times \mathfrak{X}(\mathcal{M})$ as:

$$R(\omega, X, Y, Z) := \omega(\nabla_X \nabla_Y Z - \nabla_Y \nabla_X Z - \nabla_{[X,Y]} Z) \tag{3.33}$$

Alternatively, the curvature can be viewed as a map from $\mathfrak{X}(\mathcal{M}) \times \mathfrak{X}(\mathcal{M})$ to the space of differential operators on $\mathfrak{X}(\mathcal{M})$ such that:

$$R(X,Y) = \nabla_X \nabla_Y - \nabla_Y \nabla_X - \nabla_{[X,Y]} \tag{3.34}$$

The fact that these are indeed tensors, i.e. they are linear in each argument, can be shown by direct calculation, recalling that [fX, Y] = f[X, Y] - Y(f)X.

Proposition 3.2.2. On the coordinate basis $\{\partial_{\mu}\}$ and $\{dx^{\mu}\}$ the torsion components are:

$$T^{\rho}_{\ \mu\nu} = \Gamma^{\rho}_{\mu\nu} - \Gamma^{\rho}_{\nu\mu} \tag{3.35}$$

Proof. By direct calculation:

$$T^{\rho}_{\mu\nu} = T(dx^{\rho}, \partial_{\mu}, \partial_{\nu}) = dx^{\mu}(\nabla_{\mu}\partial_{\nu} - \nabla_{\nu}\partial_{\mu} - [\partial_{\mu}, \partial_{\nu}])$$
$$= dx^{\mu}(\partial_{\mu}\partial_{\nu} - \Gamma^{\sigma}_{\mu\nu}\partial_{\sigma} - \partial_{\nu}\partial_{\mu} + \Gamma^{\sigma}_{\nu\mu}\partial_{\sigma})$$
$$= \left[\Gamma^{\sigma}_{\mu\nu} - \Gamma^{\sigma}_{\nu\mu}\right]\delta^{\mu}_{\sigma} = \Gamma^{\rho}_{\mu\nu} - \Gamma^{\rho}_{\nu\mu}$$

Interestingly, even though $\Gamma^{\rho}_{\mu\nu}$ is not a tensor, its anti-symmetric part $\Gamma^{\rho}_{[\mu\nu]} = \frac{1}{2} T^{\sigma}_{\mu\nu}$ is. Clearly, the torsion tensor is anti-symmetric in its lower indices, thus for connections which are symmetric in their lower indices the torsion is null: such connections are said to be *torsion-free*.

Proposition 3.2.3. On the coordinate basis $\{\partial_{\mu}\}$ and $\{dx^{\mu}\}$ the curvature components are:

$$R^{\sigma}_{\ \rho\mu\nu} = \partial_{\mu}\Gamma^{\sigma}_{\nu\rho} - \partial_{\nu}\Gamma^{\sigma}_{\mu\rho} + \Gamma^{\lambda}_{\nu\rho}\Gamma^{\sigma}_{\mu\lambda} - \Gamma^{\lambda}_{\mu\rho}\Gamma^{\sigma}_{\nu\lambda} \tag{3.36}$$

Proposition 3.2.4. By direct calculation:

$$\begin{split} R(dx^{\sigma},\partial_{\mu},\partial_{\nu},\partial_{\rho}) &= dx^{\sigma}(\nabla_{\mu}\nabla_{\nu}\partial_{\rho} - \nabla_{\nu}\nabla_{\mu}\partial_{\rho} - \nabla_{[\partial_{\mu},\partial_{\nu}]}\partial_{\rho}) \\ &= dx^{\sigma}(\nabla_{\mu}\nabla_{\nu}\partial_{\rho} - \nabla_{\nu}\nabla_{\mu}\partial_{\rho}) = dx^{\sigma}(\nabla_{\mu}(\Gamma^{\lambda}_{\nu\rho}\partial_{\lambda}) - \nabla_{\nu}(\Gamma^{\lambda}_{\mu\rho}\partial_{\lambda})) \\ &= dx^{\sigma}((\partial_{\mu}\Gamma^{\lambda}_{\nu\rho})\partial_{\lambda} + \Gamma^{\lambda}_{\nu\rho}\Gamma^{\tau}_{\mu\lambda}\partial_{\tau} - (\partial_{\nu}\Gamma^{\lambda}_{\mu\rho})\partial_{\lambda} - \Gamma^{\lambda}_{\mu\rho}\Gamma^{\tau}_{\nu\lambda}\partial_{\tau}) \\ &= \partial_{\mu}\Gamma^{\sigma}_{\nu\rho} - \partial_{\nu}\Gamma^{\sigma}_{\mu\rho} + \Gamma^{\lambda}_{\nu\rho}\Gamma^{\sigma}_{\mu\lambda} - \Gamma^{\lambda}_{\mu\rho}\Gamma^{\sigma}_{\nu\lambda} \end{split}$$

Clearly, the curvature tensor is anti-symmetric in its last two lower indices, i.e. $R^{\sigma}_{\ \rho\mu\nu} = R^{\sigma}_{\ \rho[\mu\nu]}$. It's also easy to show that:

$$R^{\sigma}_{\rho\mu\nu} = 2\partial_{[\mu}\Gamma^{\sigma}_{\nu]\rho} + 2\Gamma^{\sigma}_{[\mu|\lambda]}\Gamma^{\lambda}_{\nu]\rho} \tag{3.37}$$

Theorem 3.2.1. The following idendity, known as the Ricci identity, holds:

$$2\nabla_{[\mu}\nabla_{\nu]}Z^{\sigma} = R^{\sigma}_{\rho\mu\nu}Z^{\rho} - T^{\rho}_{\mu\nu}\nabla_{\rho}Z^{\sigma}$$
(3.38)

Proof. By direct calculation:

$$\begin{split} \nabla_{[\mu}\nabla_{\nu]}Z^{\sigma} &= \partial_{[\mu}(\nabla_{\nu]}Z^{\sigma}) + \Gamma^{\sigma}_{[\mu|\lambda|}\nabla_{\nu]}Z^{\lambda} - \Gamma^{\rho}_{[\mu\nu]}\nabla_{\rho}Z^{\sigma} \\ &= \partial_{[\mu}\partial_{\nu]}Z^{\sigma} + (\partial_{[\mu}\Gamma^{\sigma}_{\nu]\rho})Z^{\rho} + (\partial_{[\mu}Z^{\rho})\Gamma^{\sigma}_{\nu]\rho} + \Gamma^{\sigma}_{[\mu|\lambda|}\partial_{\nu]}Z^{\lambda} + \Gamma^{\sigma}_{[\mu|\lambda|}\Gamma^{\lambda}_{\nu]\rho}Z^{\rho} - \frac{1}{2}T^{\rho}_{\ \mu\nu}\nabla_{\rho}Z^{\sigma} \\ &= \left(\partial_{[\mu}\Gamma^{\sigma}_{\nu]\rho} + \Gamma^{\sigma}_{[\mu|\lambda|}\Gamma^{\lambda}_{\nu]\rho}\right)Z^{\rho} - \frac{1}{2}T^{\rho}_{\ \mu\nu}\nabla_{\rho}Z^{\sigma} = \frac{1}{2}R^{\sigma}_{\ \rho\mu\nu}Z^{\rho} - \frac{1}{2}T^{\sigma}_{\ \mu\nu}\nabla_{\rho}Z^{\sigma} \end{split}$$

3.2.2.2 Levi-Civita connection

The discussion on the connection has so far been independent of the metric. Starting to consider it, an important result is the fundamental theorem of Riemannian geometry.

Theorem 3.2.2 (Riemann). On a metric manifold (\mathcal{M}, g) , there exists a unique torsion-free connection that is compatible with the metric, i.e. for all $X \in \mathfrak{X}(\mathcal{M})$:

$$\nabla_X g = 0 \tag{3.39}$$

This is called the Levi-Civita connection.

Proof. WTS uniqueness: suppose such a connection exists. Then, by Leibniz:

$$X(g(Y,Z)) = \nabla_X(g(Y,Z)) = (\nabla_X g)(Y,Z) + g(\nabla_X Y,Z) + g(Y,\nabla_X,Z)$$

Since $\nabla_X g = 0$, by cyclic permutations of X, Y and Z:

$$X(g(Y,Z)) = g(\nabla_X Y, Z) + g(Y, \nabla_X Z)$$

$$Y(g(Z,X)) = g(\nabla_Y Z, X) + g(Z, \nabla_Y X)$$

$$Z(g(X,Y)) = g(\nabla_Z X, Y) + g(X, \nabla_Z Y)$$

Since the connection is torsion-free, $\nabla_X Y - \nabla_Y X = [X, Y]$, thus these equations become:

$$X(g(Y,Z)) = g(\nabla_Y X, Z) + g(\nabla_X Z, Y) + g([X,Y], Z)$$

$$Y(g(Z,X)) = g(\nabla_Z Y, X) + g(\nabla_Y X, Z) + g([Y,Z], X)$$

$$Z(g(X,Y)) = g(\nabla_X Z, Y) + g(\nabla_Z Y, X) + g([Z,X], y)$$

Adding the first two and subtracting the third:

$$g(\nabla_Y X, Z) = \frac{1}{2} [X(g(Y, Z)) + Y(g(Z, X)) + Z(g(X, Y)) - g([X, Y], Z) - g([Y, Z], X) + g([Z, X], Y)]$$

The metric is non-degenerate, thus this uniquely specifies the connection. By direct calculation it can be shown that it indeed satisfies all the properties of a connection. \Box

Proposition 3.2.5. On the coordinate basis $\{\partial_{\mu}\}$ and $\{dx^{\mu}\}$ the Levi-Civita connection's components, called *Christoffel symbols*, are:

$$\Gamma^{\lambda}_{\mu\nu} = \frac{1}{2} g^{\lambda\rho} (\partial_{\mu} g_{\nu\rho} + \partial_{\nu} g_{\mu\rho} - \partial_{\rho} g_{\mu\nu}) \tag{3.40}$$

Proof. Recalling that $[\partial_{\mu}, \partial_{\nu}] = 0$:

$$\Gamma^{\lambda}_{\mu\nu}g_{\lambda\rho} = g(\nabla_{\mu}\partial_{\nu},\partial_{\rho}) = \frac{1}{2}(\partial_{\nu}g_{\mu\rho} + \partial_{\mu}g_{\nu\rho} - \partial_{\rho}g_{\mu\nu})$$

Example 3.2.1. In flat space \mathbb{R}^n , endowed with either Euclidean or Minkowski metric, it is always possibile to choose Cartesian coordinates, in which case the Christoffel symbols vanish. Being the Riemann tensor a genuine tensor, it therefore will vanish in all possibile coordinate systems on \mathbb{R}^n , even in those with $\Gamma^{\rho}_{\mu\nu} \neq 0$: this expresses the flatness of \mathbb{R}^n .

3.2.2.3 Gauss' theorem

The divergence theorem (or Gauss' theorem) states that the integral of a total derivative is a boundary term. It is possible to express this theorem on curved manifolds in a convenient way.

Lemma 3.2.1. $\Gamma^{\mu}_{\mu\nu} = \frac{1}{\sqrt{g}} \partial_{\nu} \sqrt{g}$.

Proof. A useful identity for diagonalizable matrices: $\operatorname{tr} \log A = \log \det A$. Thus (WLOG $\det g > 0$):

$$\Gamma^{\mu}_{\mu\nu} = \frac{1}{2}g^{\mu\rho}\partial_{\nu}g_{\mu\rho} = \frac{1}{2}\operatorname{tr}(g^{-1}\partial_{\nu}g) = \frac{1}{2}\operatorname{tr}(\partial_{\nu}\log g) = \frac{1}{2}\partial_{\nu}\log\det g = \frac{1}{\sqrt{\det g}}\partial_{\nu}\sqrt{\det g}$$

Theorem 3.2.3 (Gauss). Given a Riemannian manifold (\mathcal{M},g) , consider a region $M \subseteq \mathcal{M}$ with boundary ∂M and let n^{μ} be an outward-pointing unit vector orthogonal to ∂M . Then, for any vector field X^{μ} on M:

$$\int_{M} d^{n}x \sqrt{g} \nabla_{\mu} X^{\mu} = \int_{\partial M} d^{n-1}x \sqrt{\gamma} n_{\mu} X^{\mu}$$
(3.41)

where γ_{ij} is the pull-back of the metric to ∂M and $\gamma \equiv \det \gamma_{ij}$.

Proof. From Lemma 3.2.1:

$$\sqrt{g} \nabla_{\mu} X^{\mu} = \sqrt{g} \left(\partial_{\mu} X^{\mu} + \Gamma^{\mu}_{\mu\nu} X^{\nu} \right) = \sqrt{g} \left(\partial_{\mu} X^{\mu} + X^{\nu} \frac{1}{\sqrt{g}} \partial_{\nu} \sqrt{g} \right) = \partial_{\mu} (\sqrt{g} X^{\mu})$$

The integral becomes:

$$\int_{M} d^{n}x \sqrt{g} \, \nabla_{\mu} X^{\mu} = \int_{M} d^{n}x \, \partial_{\mu} (\sqrt{g} X^{\mu})$$

This is the integral of an ordinary partial derivative, so the ordinary divergence theorem applies. To evaluate the integral on the boundary, it is convenient to pick coordinates so that ∂M is a

surface at constant x^n . Moreover, to simplify the proof, the possible metrics will be restricted to $g_{\mu\nu} = \text{diag}(\gamma_{ij}, N^2)$. By usual integration rules:

$$\int_{M} d^{n}x \, \partial_{\mu}(\sqrt{g}X^{\mu}) = \int_{\partial M} d^{n-1}x \sqrt{\gamma N^{2}} X^{n}$$

The unit normal vector is $n^{\mu} = (0, \dots, 0, \frac{1}{N})$, so that $g_{\mu\nu}n^{\mu}n^{\nu} = 1$, therefore $n_{\mu} = g_{\mu\nu}n^{\nu} = (0, \dots, 0, N)$. The proof is then concluded because:

$$\int_{\partial M} d^{n-1}x \sqrt{\gamma N^2} X^n = \int_{\partial M} d^{n-1}x \sqrt{\gamma} n_{\mu} X^{\mu}$$

Note that this theorem holds on Lorentzian manifolds too, with the condition that ∂M must be purely timelike or purely spacelike, ensuring that $\gamma \neq 0$ at any point.

3.2.2.4 Maxwell action

Consider spacetime as a manifold \mathcal{M} . The electromagnetic field can be described by a form on this manifold: indeed, the electromagnetic gauge field $A_{\mu} = (\phi, \mathbf{A})$ is to be thought as the components of a one-form $A = A_{\mu}(x)dx^{\mu}$. The exterior derivative of this form is a 2-form F = dA:

$$F = \frac{1}{2} F_{\mu\nu} dx^{\mu} \wedge dx^{\nu} = \frac{1}{2} (\partial_{\mu} A_{\nu} - \partial_{\nu} A_{\mu}) dx^{\mu} \wedge dx^{\nu}$$

The components $F_{\mu\nu}$ are in reality the components of a tensor, the Faraday tensor. By construction, a useful identity holds, sometimes called the Bianchi identity:

$$dF = 0 (3.42)$$

From this identity derive two Maxwell equations: $\nabla \cdot \mathbf{B} = 0$ and $\nabla \times \mathbf{E} + \partial_t \mathbf{B} = 0$. Moreover, note that the gauge field is not unique: the gauge transformation $A \mapsto A + d\alpha$, which equals $A_{\mu} \mapsto A_{\mu} + \partial_{\mu} \alpha$, leaves F unchanged.

To study the dynamics of these fields, an action is needed: Differential Geometry allows very few actions to be written down.

For example, suppose that on the considered manifold no metric is defined. To integrate over \mathcal{M} a 4-form is needed, but F is a 2-form, thus the only possible action is:

$$S_{\text{top}} = -\frac{1}{2} \int F \wedge F \tag{3.43}$$

The integrand become $dx^0dx^1dx^2dx^3\mathbf{E}\cdot\mathbf{B}$. Actions of this kind, independent of the metric, are called topological actions and are of no interest in classical physics: in fact, $F \wedge F = d(A \wedge F)$, so the action is a total derivative and doesno't affect the equations of motion.

To construct an action of classical interest, a metric is needed. This allows to introduce a second 2-form, $\star F$, so to construct the *Maxwell action*:

$$S_{\rm M} = -\frac{1}{2} \int F \wedge \star F \tag{3.44}$$

The integrand can then be expanded as:

$$S_{\rm M} = -\frac{1}{4} \int d^4 x \sqrt{g} \, g^{\mu\nu} g^{\rho\sigma} F_{\mu\rho} F_{\nu\sigma} = -\frac{1}{4} \int d^4 x \sqrt{g} F^{\mu\nu} F_{\mu\nu}$$

In flat spacetime $F^{\mu\nu}F_{\mu\nu}=2(\mathbf{B}^2-\mathbf{E}^2)$. In a general curved spacetime, the equation of motion resulting from the variation of the Maxwell action is $d \star F = 0$.

To complete the theory, consider a gauge field coupled to a current, described by a one-form J. The Maxwell action then becomes:

$$S_{\rm M} = \int -\frac{1}{2} F \wedge \star F + A \wedge \star J \tag{3.45}$$

This action must retain its gauge invariance, but under $A \mapsto A + d\alpha$ it transforms as $S_{\rm M} \mapsto S_{\rm M} + \int d\alpha \wedge \star J$, therefore, after integrating by parts, the condition of gauge invariance translates to:

$$d \star J = 0 \tag{3.46}$$

This is current conservation in the language of forms. Varying the action in Eq. 3.44 now leads to the Maxwell equations with source terms:

$$d \star F = \star J \tag{3.47}$$

To define electric and magnetic charges, integrate over submanifolds. Consider a three-dimensional spatial submanifold Σ : the electric charge in Σ is defined as:

$$Q_e(\Sigma) := \int_{\Sigma} \star J \tag{3.48}$$

This agrees with the usual definition in flat spacetime $Q_e = \int_{\Sigma} d^3x J^0$. Using the equations of motion and Stokes' theorem, a general form of Gauss' law is obtained:

$$Q_e(\Sigma) = \int_{\partial \Sigma} \star F \tag{3.49}$$

Similarly, the magnetic charge in Σ is defined as:

$$Q_m(\Sigma) := \int_{\partial \Sigma} F \tag{3.50}$$

The non-existence of magnetic charges, following from Bianchi identity, can be evaded in topologically interesting manifolds.

From charge conservation in Eq. 3.46, it follows that the electric charge in a region cannot change, unless current flows in or out of that region. Consider a cylindrical region of spacetime V, ending in two spatial hypersurfaces Σ_1 and Σ_2 : its boundary is $\partial V = \Sigma_1 \cup \Sigma_2 \cup B$, where B is a cylindrical timelike hypersurface. The statement that no current flows in or out of V means that $J|_B = 0$. Then:

$$Q_e(\Sigma_1) - Q_e(\Sigma_2) = \int_{\Sigma_1} \star J - \int_{\Sigma_2} \star J = \int_{\partial V} \star J - \int_B \star J = \int_{\partial V} \star J = \int_V d \star J = 0$$

Thus, electric charge in remains constant in time.

Maxwell equations from connections First note that, given the gauge field $A \in \Lambda^1(\mathcal{M})$, the field strength can be expressed via covariant derivatives:

$$F_{\mu\nu} = \partial_{\mu}A_{\nu} - \partial_{\nu}A_{\mu} = \nabla_{\mu}A_{\nu} - \nabla_{\nu}A_{\mu}$$

The Christoffel symbols cancel out due to anti-symmetry: this is what allows to define the exterior derivative without introducing connections first.

Proposition 3.2.6. Current conservation can be written as: $d \star J = 0 \Leftrightarrow \nabla_{\mu} J^{\mu} = 0$.

Proof. Recalling Lemma 3.2.1:

$$\nabla_{\mu}J^{\mu} = \partial_{\mu}J^{\mu} + \Gamma^{\mu}_{\mu\rho}J^{\rho} = \partial_{\mu}J^{\mu} + \partial_{\rho}(\log\sqrt{g})J^{\rho} = \frac{1}{\sqrt{g}}\partial_{\mu}(\sqrt{g}J^{\mu}) \propto d \star J$$

As an aside, in general the divergence in different coordinate systems can be computed using the formula $\nabla_{\mu}J^{\mu} = \frac{1}{\sqrt{g}}\partial_{\mu}(\sqrt{g}J^{\mu})$.

Proposition 3.2.7. $d \star F = \star J \Leftrightarrow \nabla_{\mu} F^{\mu\nu} = J^{\nu}$.

Proof. Recalling Lemma 3.2.1:

$$\nabla_{\mu}F^{\mu\nu} = \partial_{\mu}F^{\mu\nu} + \Gamma^{\mu}_{\mu\rho}F^{\rho\nu} + \Gamma^{\nu}_{\mu\rho}F^{\mu\rho} = \frac{1}{\sqrt{g}}\partial_{\mu}(\sqrt{g}F^{\mu\nu})$$

where $\Gamma^{\nu}_{\mu\rho}F^{\mu\rho}=0$ because $\Gamma^{\nu}_{\mu\rho}$ is symmetric in μ and ρ , while $F^{\mu\rho}$ is anti-symmetric. The proof follows recalling the definition of the Hodge dual in Eq. 3.17.

3.3 Parallel transport

The connection connects tangent spaces, or more generally any tensor vector space, at different points of the manifold: this map is called *parallel transport* and it's necessary for the definition of differentiation.

Definition 3.3.1. Consider a vector field X and some associated integral curve γ , with coordinates $x^{\mu}(\tau)$ such that:

$$X^{\mu}\big|_{\gamma} = \frac{dx^{\mu}(\tau)}{d\tau}$$

A tensor field T is said to be parallely transported along γ if:

$$\nabla_X T = 0 \tag{3.51}$$

Suppose that γ connects two points $p, q \in \mathcal{M}$: Eq. 3.51 provides a map from the tensor vector space defined at p to that defined at q. To illustrate this, consider the parallel transport of a vector field Y:

$$X^{\nu}(\partial_{\nu}Y^{\mu} + \Gamma^{\mu}_{\nu\rho}Y^{\rho}) = 0$$

Evaluating this equation on γ , considering $Y^{\mu} = Y^{\mu}(x(\tau))$:

$$\frac{dY^{\mu}}{d\tau} + X^{\nu}\Gamma^{\mu}_{\nu\rho}Y^{\rho} = 0 \tag{3.52}$$

These are a set of coupled ODEs, thus, given an initial condition (ex.: at $\tau = 0$, i.e. at p), these equations can be solved to find a unique vector at each point along the curve.

Note that the parallel transport depends both on the path (characterized by the vector field X) and on the connection.

Definition 3.3.2. Given a vector field X, a *geodesic* is a curve tangent to X such that:

$$\nabla_X X = 0 \tag{3.53}$$

Proposition 3.3.1. A geodesic is described by the geodesic equation:

$$\frac{d^2x^{\mu}}{d\tau^2} + \Gamma^{\mu}_{\nu\rho} \frac{dx^{\nu}}{d\tau} \frac{dx^{\rho}}{d\tau} = 0 \tag{3.54}$$

Proof. From the above calculations, along γ :

$$0 = \frac{dX^{\mu}}{d\tau} + \Gamma^{\mu}_{\nu\rho}X^{\nu}X^{\rho} = \frac{d^2x^{\mu}}{d\tau^2} + \Gamma^{\mu}_{\nu\rho}\frac{dx^{\nu}}{d\tau}\frac{dx^{\rho}}{d\tau}$$

For the Levi-Civita connection $\nabla_X g = 0$. If $Y \in \mathfrak{X}(\mathcal{M})$ is parallely transported along a geodesic associated to $X \in \mathfrak{X}(\mathcal{M})$, then $\nabla_X Y = \nabla_X X = 0$, therefore $\frac{d}{d\tau}g(X,Y) = 0$: this ensures that the two tangent vectors always make the same angles along the geodesic.

3.3.1 Normal coordinates

Theorem 3.3.1. Given a Riemannian manifold (\mathcal{M}, g) and $p \in \mathcal{M}$, in a neighbourhood of p it's always possible to find coordinates, called normal coordinates, such that $g_{\mu\nu}(p) = \delta_{\mu\nu}$ and $g_{\mu\nu,\rho}(p) = 0$.

Proof. By brute force, consider initial coordinates y^{μ} and find a change of coordinates x^{μ} which satisfy the requirements:

$$\frac{\partial y^{\rho}}{\partial x^{\mu}} \frac{\partial y^{\sigma}}{\partial x^{\nu}} \tilde{g}_{\rho\sigma} = g_{\mu\nu}$$

WLOG p is the origin of both sets of coordinates, so:

$$y^{\rho} = \frac{\partial y^{\rho}}{\partial x^{\mu}} \bigg|_{x=0} x^{\mu} + \frac{1}{2} \frac{\partial^{2} y^{\rho}}{\partial x^{\mu} \partial x^{\nu}} \bigg|_{x=0} x^{\mu} x^{\nu} + \dots$$

This, together with the Taylor expansion og $\tilde{g}_{\rho\sigma}$, can be inserted in the transformation equation of the metric, thus finding a set of PDEs for each power of x, which can be solved to characterize the normal coordinates. For example, the first condition is:

$$\left. \frac{\partial y^{\rho}}{\partial x^{\mu}} \right|_{x=0} \frac{\partial y^{\sigma}}{\partial x^{\nu}} \right|_{x=0} \tilde{g}_{\rho\sigma}(p) = \delta_{\mu\nu}$$

Given any $\tilde{g}_{\rho\sigma}(p)$, it's always possible to find $\frac{\partial y}{\partial x}$ that satisfies this condition. In fact, if $\dim_{\mathbb{R}} \mathcal{M} = n$, the Jacobian of the transformation has n^2 independent elements and the equation above puts $\frac{1}{2}n(n+1)$ constraints: the remaining free parameters are $\frac{1}{2}n(n-1)$, which is precisely the dimension of SO(n), the symmetry group of the flat metric. A similar counting shows that $g_{\mu\nu,\rho}(p) = 0$ puts $\frac{1}{2}n^2(n+1)$ constraints, precisely the number of independent elements of the Hessian of the transformation. \square

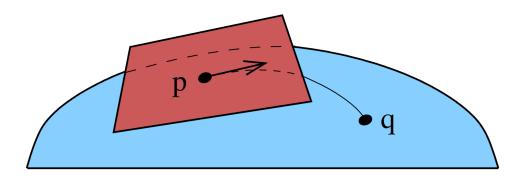


Figure 3.1: Visualization of the exponential map.

This theorem holds for Lorentzian manifolds too, but the flat metric is now $\eta_{\mu\nu}$ and its symmetry group is SO(1, n-1). The condition $g_{\mu\nu,\rho}(p)=0$ implies that $\Gamma^{\mu}_{\nu\rho}(p)=0$, but generally the Christoffel symbols won't vanish away from p. Note, however, that it's not generally possible to ensure the vanishing of second derivatives too: indeed, $g_{\mu\nu,\rho\sigma}(p)=0$ would put $\frac{1}{4}n^2(n+1)^2$ constraints, but the independent $\frac{\partial^3 y}{\partial^3 x}$ terms are $\frac{1}{6}n^2(n+1)(n+2)$, thus leaving $\frac{1}{12}n^2(n^2-1)$ free terms: this is precisely the number of independent components of the Riemann tensor, therefore in general it's not possible to pick coordinates as to make the Riemann tensor vanish too.

3.3.1.1 Exponential map

A simple way to construct normal coordinates is the following: given a tangent vector $X_p \in T_p \mathcal{M}$, there is a unique affinely parametrized geodesic through p with tangent vector X_p at p; then, any point q in the neighbourhood of p is labelled by the coordinates of the geodesic that takes from p to q a fixed amount of time.

Analytically, introducing a coordinate system (not necessarily normal) \tilde{x}^{μ} in the neighbourhood of p, an affinely parametrized geodesic solves Eq. 3.54, with initial conditions $\frac{\partial \tilde{x}^{\mu}}{\partial \tau}|_{\tau=0} = \tilde{X}^{\mu}_{p}$ and $\tilde{x}^{\mu}(\tau=0)=0$ that make the solution unique. The uniqueness of the solution allows to define a map $\text{Exp}: T_{p}\mathcal{M} \to \mathcal{M}$, called the exponential map, which acts as follows: given $X_{p} \in T_{p}\mathcal{M}$, construct the appropriate geodesic as above and follow it for a fixed affine distance, conventionally $\tau=1$, to get a new point $q \in \mathcal{M}$. See Fig. 3.1 for visual aid. Obviously, there may be points which cannot be reached from p by geodesics, or there may be tangent vectors X_{p} for which Exp is ill-defined: in General Relativity, this occurs when spacetime has singularities, but these are not relevant issues. Pick a basis $\{e_{\mu}\}$ of $T_{p}\mathcal{M}$. Then $\text{Exp}: T_{p}\mathcal{M} \ni X^{\mu}e_{\mu} \mapsto q \in \mathcal{M}$, thus it is possible to assign coordinates in the neighbourhood of p such that $x^{\mu}(q) = X^{\mu}$: these are the normal coordinates. To show this, note that if $\{e_{\mu}\}$ is orthonormal, then the geodesics will point in orthogonal directions, ensuring that $g_{\mu\nu}(p) = \delta_{\mu\nu}$. Now, fix a point q associated to a given tangent vector $X_{p} \in T_{p}\mathcal{M}$: this means that q is at distance $\tau=1$ from p along the given geodesic. Note that the geodesic equation is homogeneous in τ , thus in general $\text{Exp}: \tau X_{p} \mapsto x^{\mu}(\tau) = \tau X^{\mu}$, which means that geodesics take a simple form in these coordinates:

$$x^{\mu}(\tau) = \tau X^{\mu}$$

Being these geodesics, they must solve Eq. 3.54, that is:

$$\Gamma^{\mu}_{\nu\rho}(x(\tau))X^{\nu}X^{\rho} = 0$$

which holds at any point along the geodesic, i.e. at any $\tau \in \mathbb{R}^+$. At most points $x(\tau)$, this equation only holds for those choices of X^{μ} tangent to the geodesics. However, at x(0) = 0, i.e. at p, it must

hold for any tangent vector: this means that $\Gamma^{\mu}_{(\nu\rho)}(p) = 0$ which, for a torsion-free connection, ensures that $\Gamma^{\mu}_{\nu\rho}(p) = 0$. But vanishing Christoffel symbols imply a vanishing first derivative of the metric: for the Levi-Civita connection $2g_{\mu\sigma}\Gamma^{\sigma}_{\nu\rho} = g_{\mu\nu,\rho} + g_{\mu\rho,\nu} - g_{\nu\rho,\mu}$, thus symmetrizing $(\mu\nu)$ cancels the last two terms, leaving an identity that, evaluated at p, gives $g_{\mu\nu,\rho}(p) = 0$. Hence, these are indeed normal coordinates.

3.3.1.2 Equivalence principle

Normal coordinates are conceptually important in General Relativity: an observer at point p who parametrizes their immediate surroundings using coordinates constructed by geodesics will experience a locally flat metric. This is precisely Einstein's equivalence principle: any free-falling observer, performing local experiments, will not experience a gravitational field. The formal definition of free-falling observer is an observer which follows geodesics, while the local lack of gravitational field means $g_{\mu\mu}(p) = \eta_{\mu\nu}$. In this context, normal coordinates are called *local inertial frame*.

To understand what "local" means, note that there is a way to distinguish whether a gravitational field is present at p: a non-vanishing Riemann tensor. This depends on the second derivatives of the metric, which in general will be non-vanishing. However, to measure the effects of the Riemann tensor, one needs to compare the results of experiments at p and at a nearby point q: this is a non-local observation.

3.3.2 Curvature and torsion

With reference to Fig. 3.2, consider a tangent vector $Z_p \in T_p \mathcal{M}$ and two vector fields $X, Y \in \mathfrak{X}(\mathcal{M})$: [X,Y]=0, i.e. they are linearly independent. Construct two curved γ, γ' as in figure, both leading to a point $r \in \mathcal{M}$ which, for simplicity, is close to p. It is possible to impose normal coordinates centered at p such that $x^{\mu}=(\tau,\sigma,\ldots)$, so that $X=\frac{\partial}{\partial \tau}$ and $Y=\frac{\partial}{\partial \sigma}$: then $x^{\mu}(p)=(0,0,0,\ldots)$, $X^{\mu}(q)=(\delta\tau,0,0,\ldots)$, $x^{\mu}(s)=(0,\delta\sigma,0,\ldots)$ and $x^{\mu}(r)=(\delta\tau,\delta\sigma,0,\ldots)$, with $\delta\tau,\delta\sigma$ small.

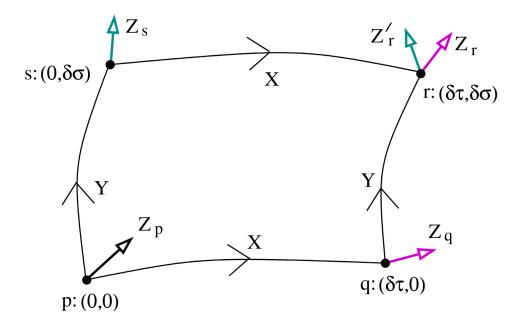


Figure 3.2: Parallel transport along different paths.

First, parallel transport Z_p along X to Z_q , so that Z^{μ} solves:

$$\frac{dZ^{\mu}}{d\tau} + X^{\nu} \Gamma^{\mu}_{\nu\rho} Z^{\rho} = 0$$

In normal coordinates $\Gamma^{\mu}_{\nu\rho}(p) = 0$, thus $\frac{dZ^{\mu}}{d\tau}\Big|_{\tau=0} = 0$ and the Taylor expansion is:

$$\begin{split} Z_q^\mu &= Z_p^\mu + \frac{\delta \tau^2}{2} \frac{d^2 Z^\mu}{d\tau^2} \bigg|_{\tau=0} + o(\delta \tau^3) \\ &= Z_p^\mu - \frac{\delta \tau^2}{s} \left[X^\nu Z^\rho \frac{d\Gamma^\mu_{\nu\rho}}{d\tau} + \frac{dX^\nu}{d\tau} Z^\rho \Gamma^\mu_{\nu\rho} + X^\nu \frac{dZ^\rho}{d\tau} \Gamma^\mu_{\nu\rho} \right]_{\tau=0} + o(\delta \tau^3) \\ &= Z_p^\mu - \frac{\delta \tau^2}{2} X^\nu Z^\rho \frac{d\Gamma^\mu_{\nu\rho}}{d\tau} \bigg|_{\tau=0} + o(\delta \tau^3) = Z_p^\mu - \frac{\delta \tau^2}{2} \left[X^\nu X^\sigma Z^\rho \Gamma^\mu_{\nu\rho,\sigma} \right]_p + o(\delta \tau^3) \end{split}$$

where $\frac{d}{d\tau} = X^{\sigma} \partial_{\sigma}$. Now, Z_q needs to be parallely transported along Y to Z_r , but this time $\frac{dZ^{\mu}}{d\sigma}|_{\sigma=0}$ doesn't vanish, in general. From the parallel transport equation:

$$\frac{dZ^{\mu}}{d\sigma}\bigg|_{\sigma=0} = -\left[Y^{\nu}Z^{\rho}\Gamma^{\mu}_{\nu\rho}\right]_{q} = -\left[Y^{\nu}Z^{\rho}X^{\sigma}\Gamma^{\mu}_{\nu\rho,\sigma}\right]_{p}\delta\tau + o(\delta\tau^{2})$$

The expansions of Y^{ν} and Z^{ρ} at leading order multiply $\Gamma^{\mu}_{\nu\rho}(p) = 0$, thus only contribute to higher order terms. Next order in $\delta\sigma$:

$$\left.\frac{d^2Z^\mu}{d\sigma^2}\right|_{\sigma=0} = -\left[\left(\frac{dY^\nu}{d\sigma}Z^\rho + Y^\nu\frac{dZ^\rho}{d\sigma}\right)\Gamma^\mu_{\nu\rho} + Y^\nu Z^\rho\frac{d\Gamma^\mu_{\nu\rho}}{d\sigma}\right]_q = -\left[Y^\nu Y^\sigma Z^\rho\Gamma^\mu_{\nu\rho,\sigma}\right]_p + o(\delta\tau)$$

The complete expansion thus is:

$$Z_r^{\mu} = Z_q^{\mu} - \left[Y^{\nu} Z^{\rho} X^{\sigma} \Gamma_{\nu\rho,\sigma}^{\mu} \right]_p \delta \tau \delta \sigma - \frac{1}{2} \left[Y^{\nu} Y^{\sigma} Z^{\rho} \Gamma_{\nu\rho,\sigma}^{\mu} \right]_p \delta \sigma^2 + o(\delta^3)$$

$$= Z_p^{\mu} - \frac{1}{2} \Gamma_{\nu\rho,\sigma}^{\mu}(p) \left[X^{\nu} X^{\sigma} Z^{\rho} \delta \tau^2 + 2 Y^{\nu} Z^{\rho} X^{\sigma} \delta \tau \delta \sigma + Y^{\nu} Y^{\sigma} Z^{\rho} \delta \sigma^2 \right] + o(\delta^3)$$

Parallel transport along γ' leads to a similar expression (exchange of X and Y):

$$Z_r^{\prime\mu} = Z_p^{\mu} - \frac{1}{2} \Gamma_{\nu\rho,\sigma}^{\mu}(p) \left[Y^{\nu} Y^{\sigma} Z^{\rho} \delta \sigma^2 + 2 X^{\nu} Z^{\rho} Y^{\sigma} \delta \sigma \delta \tau + X^{\nu} X^{\sigma} Z^{\rho} \delta \tau^2 \right] + o(\delta^3)$$

The difference between the parallely transported tangent vectors to leading order is:

$$\Delta Z_r^\mu = Z_r^\mu - {Z_r'}^\mu = -\left[\Gamma^\mu_{\nu\rho,\sigma} - \Gamma^\mu_{\sigma\rho,\nu}\right]_p \left[Y^\nu Z^\rho X^\sigma\right]_p \delta\tau \delta\sigma + o(\delta^3)$$

Recalling that $\Gamma^{\mu}_{\nu\rho}(p) = 0$, it is possible to write:

$$\Delta Z_r^{\mu} = -\left[R^{\mu}_{\rho\sigma\nu} Y^{\nu} Z^{\rho} X^{\sigma}\right]_p \delta\sigma\delta\tau + o(\delta^3)$$
(3.55)

It would be possible to evaluate the expression at r too, as it would differ only by higher order terms. Although the calculation was carried in a particular choice of coordinates, Eq. 3.55 is a tensor relation, therefore it must hold in all coordinate systems: the Riemann tensor thus determines the path dependence of parallel transport.

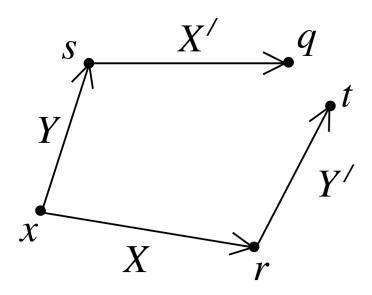


Figure 3.3: Visualization of torsion.

3.3.2.1 Torsion

Consider two tangent vectors $X_p, Y_p \in T_p \mathcal{M}$ and a coordinate system x^{μ} such that $X_p = X^{\mu} \partial_{\mu}$ and $Y_p = Y^{\mu} \partial_{\mu}$. If $p: x^{\mu}$, as in Fig. 3.3 construct $r, s \in \mathcal{M}$ such that $r: x^{\mu} + \varepsilon X^{\mu}$ and $s: x^{\mu} + \varepsilon Y^{\mu}$, with ε an infinitesimal parameter. Now, parallel transport $X_p \in T_p \mathcal{M}$ along the direction of Y_p to $X'_s \in T_s \mathcal{M}$ and $Y_p \in T_p \mathcal{M}$ along X_p to $Y'_r \in T_r \mathcal{M}$; their components will be:

$$X'_{s} = \left(X^{\mu} - \varepsilon \Gamma^{\mu}_{\nu\rho} Y^{\nu} X^{\rho}\right) \qquad Y'_{r} = \left(Y^{\mu} - \varepsilon \Gamma^{\mu}_{\nu\rho} X^{\nu} Y^{\rho}\right)$$

Repeating this process, starting from point s and moving along the direction of X'_s , a new point $q \in \mathcal{M}$ is determined, with coordinates:

$$q: x^{\mu} + \varepsilon \left(X^{\mu} + Y^{\mu} \right) - \varepsilon^2 \Gamma^{\mu}_{\nu\rho} Y^{\nu} X^{\rho}$$

Analogously, starting at point r and moving along the direction of Y'_r , a new point $t \in \mathcal{M}$ is determined, with coordinates:

$$t: x^{\mu} + \varepsilon \left(X^{\mu} + Y^{\mu} \right) - \varepsilon^{2} \Gamma^{\mu}_{\nu\rho} X^{\nu} Y^{\rho}$$

If the connection is torsion-free, then $q \equiv t$. On the other hand, if $T^{\mu}_{\nu\rho} \neq 0$, the parallelogram fails to close, as in Fig. 3.3.

3.3.3 Geodesic deviation

Definition 3.3.3. Given a one-parameter family of geodesics $\{x^{\mu}(\tau;s)\}_{s\in\mathbb{R}}$ on a manifold \mathcal{M} , the tangent vector field and the *deviation vector* field are defined as:

$$X^{\mu} := \frac{\partial x^{\mu}}{\partial \tau} \bigg|_{s} \qquad S^{\mu} := \frac{\partial x^{\mu}}{\partial s} \bigg|_{\tau}$$
 (3.56)

The meaning of these vector fields is evident: the tangent vector field fixes a particular geodesics (i.e. a particular s) and assigns at each point of the geodesic its tangent vector, while the deviation

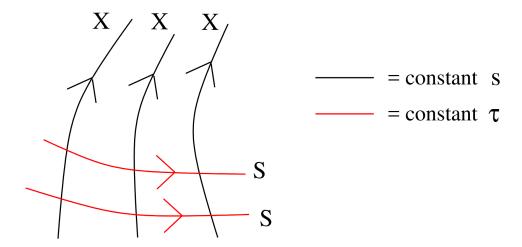


Figure 3.4: A one-parameter family of geodesics generated by X.

vector field fixes a particular value of the affine parameter τ and assigns at each point with this value a vector which takes to a nearby geodesic (at the same τ).

The family of geodesics sweeps a surface embedded in the manifold, so there is freedom in the choice of coordinates s and τ . In particular, it's always possible to pick them so that $X = \frac{\partial}{\partial \tau}$ and $S = \frac{\partial}{\partial s}$, in order for them to be linearly independent: [X, S] = 0, as in Fig. 3.4. Consider a torsion-free connection, so that:

$$\nabla_X S - \nabla_S X = [X, S] = 0$$
 \Rightarrow $\nabla_X \nabla_X S = \nabla_X \nabla_S X = \nabla_S \nabla_X X + R(X, S) X$

But X is tangent to geodesics, so from Eq. 3.53 $\nabla_X X = 0$ and:

$$\nabla_X \nabla_X S = R(X, S)X \tag{3.57}$$

Restricting to an integral curve γ of the vector field X, i.e. $X^{\mu}|_{\gamma} = \frac{dx^{\mu}}{d\tau}$, the covariant derivative along γ becomes:

$$\nabla_X \big|_{\gamma} = X^{\mu} \big|_{\gamma} \nabla_{\mu} = \frac{dx^{\mu}}{d\tau} \nabla_{\mu} \equiv \frac{D}{D\tau}$$

Hence, in index notation, the change of the deviation vector along the geodesic is expressed as:

$$\frac{DS^{\mu}}{D\tau} = R^{\mu}_{\ \nu\rho\sigma} X^{\rho} S^{\sigma} X^{\nu} \tag{3.58}$$

This can be interpreted as the relative acceleration of neighbouring geodesics, and it is determined by the Riemann tensor. Experimentally, such geodesic deviations are called *tidal forces*.

3.4 Riemann tensor

Recall Eq. 3.36 for the components of the Riemann tensor $R^{\sigma}_{\rho\mu\nu}$: it is manifestly anti-symmetric in its last two indices, but there are also other subtle symmetries when using the Levi-Civita connection.

Proposition 3.4.1. On a metric manifold with a Levi-Civita connection:

$$R_{\sigma\rho\mu\nu} = -R_{\sigma\rho\nu\mu} = -R_{\rho\sigma\mu\nu} = R_{\mu\nu\sigma\rho} \tag{3.59}$$

$$R_{\sigma[\rho\mu\nu]} = 0 \tag{3.60}$$

Proof. Set normal coordinates centered at a point p: then, $\Gamma^{\mu}_{\nu\rho} = 0$ and $\partial_{\mu}g^{\lambda\sigma} = 0$ at that point. At p, the Riemann tensor can be written as:

$$\begin{split} R_{\sigma\rho\mu\nu} &= g_{\sigma\lambda} R^{\lambda}_{\ \rho\mu\nu} = g_{\sigma\lambda} \left[\partial_{\mu} \Gamma^{\lambda}_{\nu\rho} - \partial_{\nu} \Gamma^{\lambda}_{\mu\rho} \right] \\ &= \frac{1}{2} \left[\partial_{\mu} \left(\partial_{\nu} g_{\sigma\rho} + \partial_{\rho} g_{\nu\sigma} - \partial_{\sigma} g_{\nu\rho} \right) - \partial_{\nu} \left(\partial_{\mu} g_{\sigma\rho} + \partial_{\rho} g_{\mu\sigma} - \partial_{\sigma} g_{\mu\rho} \right) \right] \\ &= \frac{1}{2} \left[\partial_{\mu} \partial_{\rho} g_{\nu\sigma} - \partial_{\mu} \partial_{\sigma} g_{\nu\rho} - \partial_{\nu} \partial_{\rho} g_{\mu\sigma} + \partial_{\nu} \partial_{\sigma} g_{\mu\rho} \right] \end{split}$$

The symmetries are then manifest, and being these tensor equations they are valid in all coordinate systems. \Box

An important computation tool is the *Bianchi identity*.

Theorem 3.4.1 (Bianchi). On a metric manifold with a Levi-Civita connection:

$$\nabla_{[\lambda} R_{\sigma\rho]\mu\nu} = 0 \qquad \Leftrightarrow \qquad R^{\sigma}_{\rho[\mu\nu;\lambda]} = 0 \tag{3.61}$$

Proof. The two equations are equivalent, so the proof is of the first one. In normal coordinates $\nabla_{\mu} = \partial_{\mu}$ at p, so schematically: $R = \partial \Gamma + \Gamma \Gamma$, thus $\nabla R = \partial R = \partial^{2} \Gamma + \Gamma \partial \Gamma = \partial^{2} \Gamma$. Explicitly:

$$\partial_{\lambda} R_{\sigma\rho\mu\nu} = \frac{1}{2} \partial_{\lambda} \left[\partial_{\mu} \partial_{\rho} g_{\nu\sigma} - \partial_{\mu} \partial_{\sigma} g_{\nu\rho} - \partial_{\nu} \partial_{\rho} g_{\mu\sigma} + \partial_{\nu} \partial_{\sigma} g_{\mu\rho} \right]$$

Anti-symmetrizing the appropriate indices yields the result.

Note that Eq. 3.60-3.61 do not require that the connection is a Levi-Civita connection, but are valid for general torsion-free connections.

3.4.1 Ricci and Einstein tensors

Definition 3.4.1. On a metric manifold, the *Ricci tensor* is defined as:

$$R_{\mu\nu} := R^{\rho}_{\ \mu\rho\nu} \tag{3.62}$$

The *Ricci scalar* is defined as:

$$R := q^{\mu\nu} R_{\mu\nu} \tag{3.63}$$

Proposition 3.4.2. On a metric manifold with a Levi-Civita connection:

$$R_{\mu\nu} = R_{\nu\mu} \tag{3.64}$$

Proof. Using Eq. 3.59:
$$R_{\mu\nu} = g^{\sigma\rho}R_{\sigma\mu\rho\nu} = g^{\rho\sigma}R_{\rho\nu\sigma\mu} = R_{\nu\mu}$$
.

Proposition 3.4.3. On a metric manifold:

$$\nabla^{\mu}R_{\mu\nu} = \frac{1}{2}\nabla_{\nu}R\tag{3.65}$$

Proof. Writing explicitly Bianchi identity:

$$\nabla_{\lambda} R_{\sigma\rho\mu\nu} + \nabla_{\sigma} R_{\rho\lambda\mu\nu} + \nabla_{\rho} R_{\lambda\sigma\mu\nu} = 0$$

Contracting with $q^{\mu\lambda}q^{\rho\nu}$:

$$\nabla^{\mu}R_{\mu\sigma} - \nabla_{\sigma}R + \nabla^{\nu}R_{\nu\sigma} = 0$$

which yields the thesis.

Definition 3.4.2. On a metric manifold, the *Einstein tensor* is defined as:

$$G_{\mu\nu} := R_{\mu\nu} - \frac{1}{2} R g_{\mu\nu} \tag{3.66}$$

Proposition 3.4.4. On a metric manifold, the Einstein tensor is covariantly constant:

$$\nabla^{\mu}G_{\mu\nu} = 0 \tag{3.67}$$

Proof. Trivial from Eq. 3.65.

3.4.2 Connection and curvature forms

This section is based on a Lorentzian manifold, but the discussion is equivalent on a Riemannian one: it's enough to swap η_{ab} with δ_{ab} as the flat metric.

3.4.2.1 Vielbeins

Although typically calculations are carried on a coordinate basis $\{e_{\mu}\}=\{\partial_{\mu}\}$, there are possible basis without such an interpretation. For example, a linear combination of a coordinate basis won't in general be a coordinate basis itself:

$$\hat{e}_a = e_a^{\ \mu} \partial_\mu \tag{3.68}$$

A particularly useful non-coordinate basis is one such that:

$$g(\hat{e}_a, \hat{e}_b) = g_{\mu\nu} e_a^{\ \mu} e_b^{\ \nu} = \eta_{ab} \tag{3.69}$$

The components $e_a^{\ \mu}$ are called *vielbeins* or *tetrads*. In this non-coordinate system, the manifold looks flat (or, at least, its patch covered by the given chart). In the following computations, greek indices are raised/lowered by the metric $g_{\mu\nu}$, while latin indices by the metric η_{ab} .

The vielbeins aren't unique, for given a set of vielbeins $e_a^{\ \mu}$ it is always possible to find a new one:

$$\tilde{e}_a^{\ \mu} = e_b^{\ \mu} (\Lambda^{-1})_a^b \tag{3.70}$$

The transformation matrix must satisfy the condition imposed by Eq. 3.69, i.e.:

$$\Lambda_a{}^c \Lambda_b{}^d \eta_{cd} = \eta_{ab} \tag{3.71}$$

These are local Lorentz transformations, because the condition is that of Lorentz transformation, but Λ is now allowed to vary over the manifold. The dual basis of one-forms $\{\hat{\theta}^a\}$ is defined by $\hat{\theta}^a(\hat{e}_b) = \delta^a_b$. The relation to the coordinate basis is:

$$\hat{\theta}^a = e^a_{\ \mu} dx^\mu \tag{3.72}$$

where the coefficients satisfy:

$$e^a_{\mu}e_b^{\mu} = \delta^a_b \tag{3.73}$$

$$e^a_{\ \mu}e^a_a{}^\nu = \delta^\nu_\mu \tag{3.74}$$

The metric is a tensor, so $g = g_{\mu\nu}dx^{\mu} \otimes dx^{\nu} = \eta_{ab}\hat{\theta}^{a} \otimes \hat{\theta}^{b}$, thus it is related to the vielbeins by:

$$g_{\mu\nu} = e^a_{\ \mu} e^b_{\ \nu} \eta_{ab} \tag{3.75}$$

3.4.2.2 Connection 1-form

On a non-coordinate basis $\{\hat{e}_a\}$, connection components are computed in the usual way:

$$\nabla_{\hat{e}_c}\hat{e}_b = \Gamma_{cb}^a \hat{e}_a \tag{3.76}$$

However, these are not the same components $\Gamma^{\mu}_{\nu\rho}$ as in a coordinate basis.

Definition 3.4.3. On a metric manifold with vielbeins, the *connection 1-form* is defined as:

$$\omega_b^a := \Gamma_{cb}^a \hat{\theta}^c \tag{3.77}$$

Note that these are really n^2 1-forms, according to values of $a, b = 1, ..., n \equiv \dim_{\mathbb{R}} \mathcal{M}$. This is also known as the *spin connection*, due to its relationship to spinors in curved spacetime.

Proposition 3.4.5. Given a local Lorentzian transformation Λ :

$$\tilde{\omega}^{a}_{b} = \Lambda^{a}_{c} \omega^{c}_{d} (\Lambda^{-1})^{d}_{b} + \Lambda^{a}_{c} (d\Lambda^{-1})^{c}_{b}$$
(3.78)

The second term reflects the second term in Prop. 3.2.1, which involves the derivative of the coordinate transformation. Important results for connection 1-forms is the *Cartan structure equation*.

Theorem 3.4.2 (Cartan). On a metric manifold with a torsion-free connection:

$$d\hat{\theta}^a + \omega^a_b \wedge \hat{\theta}^b = 0 \tag{3.79}$$

Proof. From Eq. 3.76 $\Gamma^a_{cb} = e^a_{\ \rho} e_c^{\ \mu} \nabla_{\mu} e_b^{\ \rho}$, thus, remembering $\Gamma^{\rho}_{[\mu\nu]} = 0$ (torsion-free):

$$\omega^{a}_{b} \wedge \hat{\theta}^{b} = \Gamma^{a}_{cb} \left(e^{c}_{\mu} dx^{\mu} \right) \wedge \left(e^{b}_{\nu} dx^{\nu} \right)$$

$$= e^{a}_{\rho} e^{\mu}_{c} \left(\partial_{\mu} e^{\rho}_{b} + e^{\nu}_{b} \Gamma^{\rho}_{\mu\nu} \right) \left(e^{c}_{\mu} dx^{\mu} \right) \wedge \left(e^{b}_{\nu} dx^{\nu} \right)$$

$$= e^{a}_{\rho} \underbrace{e^{\lambda}_{c} e^{c}_{\mu}}_{\delta^{\mu}_{u}} e^{b}_{\nu} \left(\partial_{\lambda} e^{\rho}_{b} + e^{\sigma}_{b} \Gamma^{\rho}_{\lambda\sigma} \right) dx^{\mu} \wedge dx^{\nu} = e^{a}_{\rho} e^{b}_{\nu} \partial_{\mu} e^{\rho}_{b} dx^{\mu} \wedge dx^{\nu}$$

But $e^b_{\ \nu}e^{\ \rho}_b = \delta^\rho_{\nu}$, so $e^b_{\ \nu}\partial_{\mu}e^{\ \rho}_b = -e^{\ \rho}_b\partial_{\mu}e^n_{\ \nu}$, hence:

$$\omega^a_{\ b}\wedge\hat\theta^b=-e^a_{\ \rho}e_b^{\ \rho}\partial_\mu e^b_{\ \nu}dx^\mu\wedge dx^\nu=-\partial_\mu e^a_{\ \nu}dx^\mu\wedge dx^\nu=-d\hat\theta^a$$

For a Levi-Civita connection, a stronger result holds.

Proposition 3.4.6. On a metric manifold with a Levi-Civita connection:

$$\omega_{ab} = -\omega_{ba} \tag{3.80}$$

Proof. Being the Levi-Civita connection compatible with the metric:

$$\begin{split} \Gamma_{abc} &= \eta_{ad} e^d_{\ \rho} e_b^{\ \mu} \nabla_{\mu} e_c^{\ \rho} = -\eta_{ad} e_c^{\ \rho} e_b^{\ \mu} \nabla_{\mu} e^d_{\ \rho} = -\eta_{cf} e^f_{\ \sigma} e_b^{\ \mu} \nabla_{\mu} (\eta_{ad} g^{\rho \sigma} e^d_{\ \rho}) \\ &= -\eta_{cf} e^f_{\ \rho} e_b^{\ \mu} \nabla_{\mu} e_a^{\ \rho} = -\Gamma_{cba} \end{split}$$

From Eq. 3.77 $\omega_{ab} = \Gamma_{acb} \hat{\theta}^c$, thus completing the proof.

Eq. 3.79-3.80 allow to quickly compute the spin connection, as they uniquely define it. Indeed, ω_{ab} being anti-symmetric means that there are $\frac{1}{2}n(n-1)$ independent 1-forms, i.e. $\frac{1}{2}n^2(n-1)$ independent components. The Cartan structure equation relates two 2-forms, each with $\frac{1}{2}n(n-1)$ independent components, thus posing $\frac{1}{2}n^2(n-1)$ constraints (as there are n equations) and uniquely fixing the spin connection.

3.4.2.3 Curvature 2-form

Recall Eq. 3.59, which holds for Levi-Civita connections. Computing the Riemann tensor in the non-coordinate basis $R^a_{bcd} = R(\hat{\theta}^a, \hat{e}_b, \hat{e}_c, \hat{e}_d)$, the anti-symmetry of the last two indices persists: $R^a_{bcd} = -R^a_{bdc}$.

Definition 3.4.4. On a metric manifold with a Levi-Civita connection, the *curvature 2-form* is defined as:

$$\mathcal{R}^{a}_{b} := \frac{1}{2} R^{a}_{bcd} \hat{\theta}^{c} \wedge \hat{\theta}^{d} \tag{3.81}$$

Again, these are really n^2 2-forms. A second Cartan structure relation holds.

Theorem 3.4.3 (Cartan). On a metric manifold with a Levi-Civita connection:

$$\mathcal{R}^{a}_{b} = d\omega^{a}_{b} + \omega^{a}_{c} \wedge \omega^{c}_{b} \tag{3.82}$$

Connection and curvature forms make computing the Riemann tensor less tedious, as exterior derivatives take significant less effort than covariant derivatives.

Part III General Relativity

Einstein Field Equations

The force of gravity is mediated by a gravitational field, which is identified with a metric $g_{\mu\nu}(x)$ on a 4d Lorentzian manifold called spacetime. This metric is a dynamical object, as all other fields in Nature, thus the laws governing its dynamics must be studied.

4.1 Einstein-Hilbert action

Differential Geometry places strict limits to the possible actions that can be formulated, ensuring it is something intrinsic to the metric and independent on the particular choice of coordinates. Spacetime is a Lorentzian manifold \mathcal{M} , thus, recalling the canonical volume form in Eq. 3.13, there need to be metric-dependent scalar function on \mathcal{M} : the obvious non-trivial choice is the Ricci scalar. The resulting action is the *Einstein-Hilbert action*:

$$S = \int d^4x \sqrt{-g}R \tag{4.1}$$

where the negative sign makes it manifest that the metric has signature (-,+,+,+). Schematically, the Riemann tensor is $R \sim \partial \Gamma + \Gamma \Gamma$, while $\Gamma \sim \partial g$, thus the Einstein-Hilbert action is second order in derivatives of the metric, like most other actions in physics.

4.1.1 Equations of motion

To determine the Euler-Lagrange equations of the Einstein-Hilbert action, consider a fixed initial metric $g_{\mu\nu}(x)$ and a perturbation $g_{\mu\nu}(x) \mapsto g_{\mu\nu}(x) + \delta g_{\mu\nu}(x)$. For the inverse metric:

$$g_{\mu\nu}g^{\nu\rho} = \delta^{\rho}_{\mu} \quad \Rightarrow \quad \delta g_{\mu\nu}g^{\nu\rho} + g_{\mu\nu}\delta g^{\nu\rho} = 0 \quad \Rightarrow \quad \delta g^{\mu\nu} = -g^{\mu\rho}g^{\nu\sigma}\delta g_{\rho\sigma}$$

Lemma 4.1.1. The variation of $\sqrt{-g}$ is:

$$\delta\sqrt{-g} = -\frac{1}{2}\sqrt{-g}\,g_{\mu\nu}\delta g^{\mu\nu} \tag{4.2}$$

Proof. Recalling that for a diagonalizable matrix $\log \det A = \operatorname{tr} \log A$, then $(\det A)^{-1}\delta(\det A) = \operatorname{tr}(A^{-1}\delta A)$. Applying this to the metric:

$$\delta\sqrt{-g} = \frac{1}{2} \frac{1}{\sqrt{-g}} \delta(-g) = \frac{1}{2} \frac{1}{\sqrt{-g}} (-g) g^{\mu\nu} \delta g_{\mu\nu} = \frac{1}{2} \sqrt{-g} g^{\mu\nu} \delta g_{\mu\nu} = -\frac{1}{2} \sqrt{-g} g_{\mu\nu} g^{\mu\nu}$$

Lemma 4.1.2. The variation of the Christoffel symbols is:

$$\delta\Gamma^{\rho}_{\mu\nu} = \frac{1}{2}g^{\rho\sigma}(\nabla_{\mu}\delta g_{\sigma\nu} + \nabla_{\nu}\delta g_{\sigma\mu} - \nabla_{\sigma}\delta g_{\mu\nu}) \tag{4.3}$$

Proof. First note that, although $\Gamma^{\rho}_{\mu\nu}$ is not a tensor, $\delta\Gamma^{\rho}_{\mu\nu}$ is, for it is the difference of two Christoffel symbols, one computed with $g_{\mu\nu}$ and the other with $g_{\mu\nu} + \delta g_{\mu\nu}$, but the extra term in the transformation law of $\Gamma^{\rho}_{\mu\nu}$ is independent of the metric, thus cancels out.

This osservation allows to work in normal coordinates at $p \in \mathcal{M}$, so that $\partial_{\rho}g_{\mu\nu} = 0$ and $\Gamma^{\rho}_{\mu\nu} = 0$ at that point. Then, to linear order in the variation, at p:

$$\delta\Gamma^{\rho}_{\mu\nu} = \frac{1}{2}g^{\rho\sigma}(\partial_{\mu}\delta g_{\sigma\nu} + \partial_{\nu}\delta g_{\sigma\mu} - \partial_{\sigma}\delta g_{\mu\nu}) = \frac{1}{2}g^{\rho\sigma}(\nabla_{\mu}\delta g_{\sigma\nu} + \nabla_{\nu}\delta g_{\sigma\mu} - \nabla_{\sigma}\delta g_{\mu\nu})$$

This is a tensor equation, hence valid in all coordinate system and, being p arbitrary, on all \mathcal{M} . \square

Lemma 4.1.3. The variation of the Ricci tensor is:

$$\delta R_{\mu\nu} = \nabla_{\rho} \delta \Gamma^{\rho}_{\rho\nu} - \nabla_{\nu} \delta \Gamma^{\rho}_{\rho\mu} \tag{4.4}$$

Proof. Working in normal coordinates, the Riemann tensor becomes $R^{\sigma}_{\mu\rho\nu} = \partial_{\rho}\Gamma^{\sigma}_{\nu\mu} - \partial_{\nu}\Gamma^{\sigma}_{\rho\mu}$, so:

$$\delta R^{\sigma}_{\ \mu\rho\nu} = \partial_{\rho} \delta \Gamma^{\sigma}_{\nu\mu} - \partial_{\nu} \delta \Gamma^{\sigma}_{\rho\mu} = \nabla_{\rho} \delta \Gamma^{\sigma}_{\nu\mu} - \nabla_{\nu} \delta \Gamma^{\sigma}_{\rho\mu}$$

This is a tensor equation, hence valid in all coordinates systems and on all \mathcal{M} . Contracting indices σ, ρ and working to leading order yields the result.

Proposition 4.1.1. The Euler-Lagrange equations of the Einstein-Hilbert action are:

$$G_{\mu\nu} = 0 \tag{4.5}$$

Proof. Varying the Einstein-Hilbert action:

$$\delta \mathcal{S} = \delta \int d^4 x \sqrt{-g} \, g^{\mu\nu} R_{\mu\nu}$$

$$= \int d^4 x \left[\delta(\sqrt{-g}) g^{\mu\nu} R_{\mu\nu} + \sqrt{-g} (\delta g^{\mu\nu}) R_{\mu\nu} + \sqrt{-g} \, g^{\mu\nu} (\delta R_{\mu\nu}) \right]$$

$$= \int d^4 x \sqrt{-g} \left[\left(-\frac{1}{2} R g_{\mu\nu} + R_{\mu\nu} \right) \delta g^{\mu\nu} + g^{\mu\nu} \left(\nabla_{\rho} \delta \Gamma^{\rho}_{\mu\nu} - \nabla_{\nu} \Gamma^{\rho}_{\rho\mu} \right) \right]$$

$$= \int d^4 x \sqrt{-g} \left[\left(R_{\mu\nu} - \frac{1}{2} R g_{\mu\nu} \right) \delta g^{\mu\nu} + \nabla_{\mu} \left(g^{\rho\nu} \nu \Gamma^{\mu}_{\rho\nu} - g^{\mu\nu} \delta \Gamma^{\rho}_{\rho\nu} \right) \right]$$

The last term is a total derivative, hence by the divergence theorem Eq. 3.41 it yields a boundary term which can be ignored (Gibbons-Hawking boundary term). The Euler-Lagrange equations are then found imposing δS for arbitrary $\delta g_{\mu\nu}$, so recalling the definition of the Einstein tensor Eq. 3.66 the proof is completed.

These equations are called *Einstein field equations* in the absence of any matter. For this reason, they can be further simplified: by contracting with $g^{\mu\nu}$ one finds R=0, thus:

$$R_{\mu\nu} = 0 \tag{4.6}$$

4.1.1.1 Dimensional analysis

The Einstein-Hilbert action in Eq. 4.1 doesn't have the right dimensions, which will be necessary when considering the presence of matter.

If x^{μ} has dimension of length, then $g_{\mu\nu}$ is dimensionless and the Ricci scalar is $[R] = L^{-2}$. Including the integration measure $[S] = L^2$, but it must have the same dimension of $[\hbar] = ML^2T^{-1}$ (i.e. energy \times times). Thus, the action with the right dimension is:

$$S = \frac{c^3}{16\pi G} \int d^4x \sqrt{-g}R \tag{4.7}$$

In the following, natural units are adopted: $c = \hbar = 1$.

4.1.1.2 Cosmological constant

It reality, Eq. 4.1 is not the simplest action allowed by Differential Geometry: in fact, a constant term could be added to the Ricci scalar. The action then becomes:

$$S = \frac{1}{16\alpha G} \int d^4x \sqrt{-g} (R - 2\Lambda)$$
 (4.8)

 Λ is called the *cosmological constant* and has dimension $[\Lambda] = L^{-2}$. The resulting field equations are:

$$R_{\mu\nu} - \frac{1}{2}Rg_{\mu\nu} = -\Lambda g_{\mu\nu} \tag{4.9}$$

Contracting with $g^{\mu\nu}$ yields $R=4\Lambda$, thus in the absence of matter:

$$R_{\mu\nu} = \Lambda g_{\mu\nu} \tag{4.10}$$

4.1.2 Diffeomorphisms

Being the metric a symmetric $\mathbb{R}^{4,4}$ matrix, it should have $\frac{1}{2} \times 4 \times 5 = 10$ degrees of freedom. However, not all the components of $g_{\mu\nu}$ are physical. Indeed, two metrics related by a change of coordinates $x^{\mu} \mapsto \tilde{x}^{\mu}(x)$ describe the same physical spacetime, thus there is a redundancy in any given representation of the metric, which removes precisely 4 degrees of freedom, leaving with only 6 remaining degrees of freedom.

Mathematically, given a diffeomorphism $\phi: \mathcal{M} \to \mathcal{M}$, it maps all fields on \mathcal{M} to a new set of fields on \mathcal{M} : the result is physically indistinguishable from the original, describing the same spacetime but in different coordinates. Such diffeomorphisms are analogous to gauge symmetries in Yang-Mills theory.

Consider a diffeomorphism $x^{\mu} \mapsto \tilde{x}^{\mu} = x^{\mu} + \delta x^{\mu}$: this can be viewed as an active change, where points in spacetime are mapped from one another, or as a passive change, in which only the coordinates of each point are affected, but the two are equivalent. This change of coordinates can be thought as generated by an infinitesimal vector field $X^{\mu} : \delta x^{\mu} = -X^{\mu}(x)$, so that the metric transforms as:

$$\tilde{g}_{\mu\nu}(\tilde{x}) = \frac{\partial x^{\rho}}{\partial \tilde{x}^{\mu}} \frac{\partial x^{\sigma}}{\partial \tilde{x}^{\nu}} g_{\rho\sigma}(x) = \left(\delta^{\rho}_{\mu} + \partial_{\mu} X^{\rho}\right) \left(\delta^{\sigma}_{\nu} + \partial_{\nu} X^{\sigma}\right) g_{\rho\sigma}(x)$$
$$= g_{\mu\nu}(x) + g_{\mu\rho}(x) \partial_{\nu} X^{\rho} + g_{\nu\rho}(x) \partial_{\mu} X^{\rho}$$

Meanwhile, the Taylor expansion around $\tilde{x} = x + \delta x$ is:

$$\tilde{g}_{\mu\nu}(\tilde{x}) = \tilde{g}_{\mu\nu}(x + \delta x) = \tilde{g}_{\mu\nu}(x) - X^{\lambda} \partial_{\lambda} \tilde{g}_{\mu\nu}(x)$$

Comparing the metrics at the same point, it is understood that it undergoes an infinitesimal change:

$$\delta g_{\mu\nu}(x) = \tilde{g}_{\mu\nu}(x) - g_{\mu\nu}(x) = X^{\lambda} \partial_{\lambda} g_{\mu\nu}(x) + g_{\mu\rho}(x) \partial_{\nu} X^{\rho} + g_{\nu\rho}(x) \partial_{\mu} X^{\rho}$$

By Eq. 2.23, this is the Lie derivative of the metric:

$$\delta g_{\mu\nu} = (\mathcal{L}_X g)_{\mu\nu} \tag{4.11}$$

Thus, an infinitesimal diffeomorphism along $X \in \mathfrak{X}(\mathcal{M})$ makes the metric change by an infinitesimal amount given by its Lie derivative along X. This can be viewed as the leading term in the Taylor expansion along X. These equations can be rewritten in a simpler form:

$$\delta g_{\mu\nu} = X^{\rho} \partial_{\rho} g_{\mu\nu} + \partial_{\nu} X_{\mu} - X^{\rho} \partial_{\nu} g_{\mu\rho} + \partial_{\mu} X_{\nu} - X^{\rho} \partial_{\mu} g_{\nu\rho} = \partial_{\mu} X_{\nu} + \partial_{\nu} X_{\mu} + 2g_{\rho\sigma} \Gamma^{\sigma}_{\mu\nu} X^{\rho}$$

Therefore:

$$\delta g_{\mu\nu} = \nabla_{\mu} X_{\nu} + \nabla_{\nu} X_{\mu} \tag{4.12}$$

This equation can be used in the path integral. In fact, insisting that $\delta S = 0$ for any $\delta g_{\mu\nu}$ gives the equations of motion; on the other hand, those variations $\delta g_{\mu\nu}$ for which $\delta S = 0$ for any metric are the *symmetries* of the action. From Prop. 4.1.1:

$$\delta \mathcal{S} = \int d^4x \sqrt{-g} \, G^{\mu\nu} \delta g_{\mu\nu} = 2 \int d^4x \sqrt{-g} \, G^{\mu\nu} \nabla_{\mu} X_{\nu}$$

Invariance for $x^{\mu} \mapsto x^{\mu} - X^{\mu}$ means that $\delta S = 0$ for all $X \in \mathfrak{X}(\mathcal{M})$, hence, integrating by parts:

$$\nabla_{\mu}G^{\mu\nu} = 0 \tag{4.13}$$

This is Bianchi identity: although it is a result from Differential Geometry, is follows from diffeomorphism invariance of the Einstein-Hilbert action. Bianchi identity is comprised of 4 equations, which make the 10 Einstein equations not completely independent: in fact, only 6 of then are independent, the same number of independent components of the metric, thus ensuring that the field equations are not overdetermined.

4.2 Simple solutions

The Einstein equations in the absence of matter are:

$$R_{\mu\nu} = \Lambda g_{\mu\nu}$$

with $\Lambda \in \mathbb{R}$.

4.2.1 Minkowski space

The simplest case is $\Lambda = 0$. The Einstein equations then reduce to $R_{\mu\nu} = 0$, with the condition of $g_{\mu\nu}$ being non-degenerate, for it is a metric and the field equations require the existence of the inverse metric $g^{\mu\nu}$. This restriction is physically unusual: it is not a holonomic constraint on the physical degrees of freedom, but an inequality det g < 0 (with required signature (-, +, +, +)) which is not found for other fields of the Standard Model.

The simplest Ricci flat metric is Minkowski space $ds^2 = -dt^2 + d\mathbf{x}^2$, but it is not the only metric obeying $R_{\mu\nu} = 0$. Another example is Schwarzschild metric.

4.2.2 de Sitter space

Consider $\Lambda > 0$. A possible ansatz is:

$$ds^{2} = -f(r)^{2}dt^{2} + f(r)^{-2}dr^{2} + r^{2}(d\vartheta^{2} + \sin^{2}\vartheta d\varphi^{2})$$

4.2.2.1 Riemann tensor

First, one needs to compute the Ricci tensor for this metric: a simple way is using the curvature form. The non-coordinate 1-forms satisfy $d^2 = \eta_{ab} \hat{\theta}^a \otimes \hat{\theta}^b$, thus:

$$\hat{\theta}^0 = f dt \qquad \hat{\theta}^1 = f^{-1} dr \qquad \hat{\theta}^2 = r d\vartheta \qquad \hat{\theta}^3 = r \sin \vartheta d\varphi$$

$$d\hat{\theta}^0 = f' dr \wedge dt \qquad d\hat{\theta}^1 = 0 \qquad d\hat{\theta}^2 = dr \wedge d\vartheta \qquad d\hat{\theta}^3 = \sin \vartheta dr \wedge d\varphi + r \cos \vartheta d\vartheta \wedge d\varphi$$

The first Cartan structure relation Eq. 3.79, together with Eq. 3.80, allow to determine the connection 1-form. For example, the first equation is $\omega_1^0 = f'fdt = f'd\hat{\theta}^0$, and then $\omega_0^1 = \omega_{10} = -\omega_{01} = \omega_1^0$. Using the other structure equation, the only non-vanishing components of the connection 1-form are:

$$\omega_{1}^{0} = \omega_{0}^{1} = f'\hat{\theta}^{0} \qquad \omega_{1}^{2} = -\omega_{2}^{1} = \frac{f}{r}\hat{\theta}^{2} \qquad \omega_{1}^{3} = -\omega_{3}^{1} = \frac{f}{r}\hat{\theta}^{3} \qquad \omega_{2}^{3} = -\omega_{3}^{2} = \frac{\cot\theta}{r}\hat{\theta}^{3}$$

The curvature 2-form can be computed from the second Cartan structure relation Eq. 3.82. For example, $\mathcal{R}^0_{\ 1} = d\omega^0_{\ 1} + \omega^0_{\ c} \wedge \omega^c_{\ 1}$, but $\omega^0_{\ c} \wedge \omega^c_{\ 1} = \omega^0_{\ 1} \wedge \omega^1_{\ 1} = 0$, thus $\mathcal{R}^0_{\ 1} = d\omega^0_{\ 1} = ((f')^2 + f''f)dr \wedge dt$. From the curvature 2-form, the Riemann tensor can be calculated via Eq. 3.81 (recall the antisymmetries of the Riemann tensor), finding the only non-vanishing components:

$$R_{0101} = ff'' + (f')^2$$
 $R_{0202} = R_{0303} = \frac{ff'}{r}$ $R_{1212} = R_{1313} = -\frac{ff'}{r}$ $R_{2323} = \frac{1 - f^2}{r^2}$

To convert back to $x^{\mu}=(t,r,\vartheta,\varphi)$, use $R_{\mu\nu\rho\sigma}=e^a_{\ \mu}e^b_{\ \nu}e^c_{\ \rho}e^d_{\ \sigma}R_{abcd}$, which is particularly easy given that the matrices $e_a^{\ \mu}$ which define the non-coordinate 1-forms are diagonal:

$$R_{trtr} = f(r)f''(r) + f'(r)^{2} \qquad R_{t\vartheta t\vartheta} = f(r)^{3}f'(r)r \qquad R_{t\varphi t\varphi} = f(r)^{3}f'(r)r\sin^{2}\vartheta$$

$$R_{r\vartheta r\vartheta} = -\frac{f'(r)r}{f(r)} \qquad R_{r\varphi r\varphi} = -\frac{f'(r)r}{f(r)}\sin^{2}\vartheta \qquad R_{\vartheta\varphi\vartheta\varphi} = (1 - f(r)^{2})r^{2}\sin^{2}\vartheta$$

4.2.2.2 Ricci tensor

Given the Riemann tensor, it is easy to check that the Ricci tensor is diagonal with components:

$$R_{tt} = -f(r)^4 R_{rr} = f(r)^3 \left[f''(r) + \frac{2f'(r)}{r} + \frac{f'(r)^2}{f(r)} \right]$$
$$R_{\varphi\varphi} = \sin^2 \vartheta R_{\vartheta\vartheta} = (1 - f(r)^2 - 2f(r)f'(r)r)\sin^2 \vartheta$$

Imposing $R_{\mu\nu} = \Lambda g_{\mu\nu}$ determines two constraints on f(r):

$$f''(r) + \frac{2f'(r)}{r} + \frac{f'(r)^2}{f(r)} = -\frac{\Lambda}{f(r)}$$
 1 - 2f(r)f'(r)r - f(r)^2 = \Lambda r^2

A solution is:

$$f(r) = \sqrt{1 - \frac{r^2}{R^2}} \qquad R^2 \equiv \frac{3}{\Lambda}$$

This determines the metric of de Sitter space:

$$ds^{2} = -\left(1 - \frac{r^{2}}{R^{2}}\right)dt^{2} + \left(1 - \frac{r^{2}}{R^{2}}\right)^{-1}dr^{2} + r^{2}(d\vartheta^{2} + \sin^{2}\vartheta d\varphi^{2})$$
(4.14)

More precisely, this is the *static patch* of de Sitter space.

4.2.2.3 de Sitter geodesics

To interpret the metric, it's useful to study its geodesics. First, note that a non-trivial $g_{tt}(r)$ term means that a particle won't remain at rest at $r \neq 0$, but it will be pushed to smaller values of $g_{tt}(r)$, i.e. larger values of r. The action of a particle in the general f(r) metric, parametrized by its proper time σ , is:

$$S = \int d\sigma \left[-f(r)^2 \dot{t}^2 + f(r)^{-2} \dot{r}^2 + r^2 (\dot{\vartheta}^2 + \sin^2 \vartheta \varphi^2) \right]$$
 (4.15)

where $\dot{x}^{\mu} \equiv \frac{dx^{\mu}}{d\sigma}$. This Lagrangian has two ignorable degrees of freedom which lead to conserved quantities: $t(\sigma)$ and $\varphi(\sigma)$, as they appear only with time derivatives. The first one has energy as the conserved quantity, while the second has angular momentum:

$$E = -\frac{1}{2} \frac{\partial L}{\partial \dot{t}} = f(r)^2 \dot{t} \tag{4.16}$$

$$\ell = \frac{1}{2} \frac{\partial L}{\partial \dot{\varphi}} = r^2 \sin^2 \vartheta \dot{\varphi} \tag{4.17}$$

The $\frac{1}{2}$ are due to the absence of the usual factor in the kinetic term. The equations of motion from the action Eq. 4.15 should be supplemented with a constraint to distinguish whether a particle is massive or massless. For a massive particle, the trajectory must be timelike, so:

$$-f(r)^{2}\dot{t}^{2} + f(r)^{-2}\dot{r}^{2} + r^{2}(\dot{\vartheta}^{2} + \sin^{2}\vartheta\dot{\varphi}^{2}) = -1$$

WLOG consider geodesics that lie in the $\vartheta = \frac{\pi}{2}$ plane, so that $\dot{\theta} = 0$ and $\sin^2 \vartheta = 1$. Replacing (t, φ) with (E, ℓ) , the constraint becomes:

$$\dot{r}^2 + V_{\text{eff}}(r) = E^2 \qquad V_{\text{eff}}(r) = \left(1 + \frac{\ell^2}{r^2}\right) f(r)^2$$
 (4.18)

For de Sitter space:

$$V_{\text{eff}}(r) = \left(1 + \frac{\ell^2}{r^2}\right) \left(1 - \frac{r^2}{R^2}\right)$$

This effective potential is plotted in Fig. 4.1: immediately, on sees that the potential pushes the particle to larger values of r. Focusing on geodesics with $\ell=0$, the potential is a harmonic repulsor: a particle stationary at r=0 is an unstable geodesic, for if it has non-zero initial velocity it will follow the trajectory:

$$r(\sigma) = R\sqrt{E^2 - 1}\sinh\frac{\sigma}{R} \tag{4.19}$$

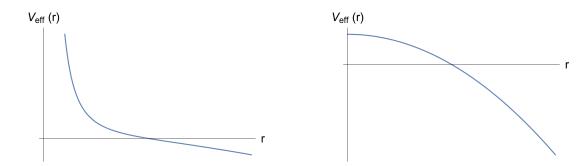


Figure 4.1: Effective potential for de Sitter geodesics, with $\ell \neq 0$ and $\ell = 0$ respectively.

Note that the metric is singular at r = R, a fact not manifest in the geodesic Eq. 4.19 which shows that any observer reaches r = R in finite proper time. Problems arise when studying the coordinate time t, which has the interpretation of the time experienced by someone stationary at r = 0; from Eq. 4.16:

$$\frac{dt}{d\sigma} = E\left(1 - \frac{r^2}{R^2}\right)^{-1}$$

This shows that $t(\sigma) \to \infty$ as $r(\sigma) \to R$: in fact, if $r(\sigma^*) = R$, then for $\sigma = \sigma^* - \varepsilon$ one has $\frac{dt}{d\sigma} = -\frac{\alpha}{\varepsilon}$, i.e. $t(\varepsilon) = -\log(\varepsilon/R)$, so $t(\varepsilon) \to \infty$ as $\varepsilon \to 0$. This means that a particle moving along the geodesic Eq. 4.19 will reach r = R in finite proper time, but a stationary observer at r = 0 will measure an infinite amount of time.

This strange behaviour is similar to what happens at the horizon of a black hole (r = 2GM): however, the Schwarzschild metric has a singularity at r = 0, while de Sitter metric looks flat at r = 0. de Sitter space seems like an inverted black hole in which particles are pushed outwars to r = R.

4.2.2.4 de Sitter embeddings

de Sitter space can be nicely embedded as a submanifold of $\mathbb{R}^{1,4}$ with metric:

$$ds^{2} = -(dX^{0})^{2} + \sum_{i=1}^{4} (dX^{i})^{2}$$
(4.20)

In particular, de Sitter metric Eq. 4.14 is a metric on the submanifold of $\mathbb{R}^{1,4}$ defined by the timelike hyperboloid:

$$-(X^{0})^{2} + \sum_{i=1}^{4} (X^{i})^{2} = R^{2}$$
(4.21)

A way of parametrizing this constraint is by imposing that $r^2 = (X^1)^2 + (X^2)^2 + (X^3)^2$, so that:

$$R^2 - r^2 = (X^4)^2 - (X^0)^2$$

The solutions are parametrized as:

$$X^{0} = \sqrt{R^{2} - r^{2}} \sinh \frac{t}{R} \qquad X^{4} = \sqrt{R^{2} - r^{2}} \cosh \frac{t}{R}$$

Computing the respective variations, along with $\sum_{i=1}^{3} (dX^i)^2 = dr^2 + r^2 d\Omega_2^2$, where $d\Omega_n^2$ is the metric element on \mathbb{S}^n , allows to show that the pull-back of the Minkowski metric Eq. 4.20 on the hyperboloid

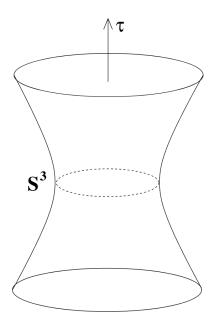


Figure 4.2: de Sitter space visualization with global coordinates.

Eq. 4.21 so parametrized gives the de Sitter metric Eq. 4.14 in the static patch coordinates. The coordinates $\{X^i\}_{i=0,\dots,4}$ so defined are not the most intuitive: they single out X^4 as special, when the constraint does no such thing, and they do not cover the whole hyperboloid, as they are limited to $X^4 \geq 0$. A better choice of parametrization is:

$$X^{0} = R \sinh \frac{\tau}{T}$$
 $X^{i} = y^{i}R \cosh \frac{\tau}{R} : \sum_{i=1}^{4} (y^{i})^{2} = 1$

Given this constraint, $\{y^i\}_{i=1,2,3,4}$ parametrize \mathbb{S}^3 . These coordinates retain more of the symmetry of de Sitter space and cover the whole space, thus are a better parametrization. The pull-back of Minkowski metric, however, gives a rather different metric on de Sitter space:

$$ds^2 = -d\tau^2 + R^2 \cosh^2 \frac{\tau}{R} d\Omega_3^2 \tag{4.22}$$

These are known as *global coordinates*, as they cover the whole space (except for singularities at the poles for any choice of coordinates on \mathbb{S}^3). Since this metric is related to that in Eq. 4.14 by a change of coordinates, it too must obey the Einstein equations. Global coordinates also show that the singularity which happens at X^4 , i.e. at r = R, is nothing but a coordinate singularity.

These coordinates provide a clearer intuition for the physics of de Sitter space: it is a time-dependent solution of the field equations in which a spatial \mathbb{S}^3 first shrinks to a minimal radius R and then expands, as shown in Fig. 4.2. The expansionary phase is a good approximation to our current universe on large scales.

4.2.3 Anti-de Sitter space

Consider $\Lambda < 0$. With the same ansatz as for de Sitter space, it's easy to find the metric for *anti-de Sitter space* (AdS):

$$ds^{2} = -\left(1 + \frac{r^{2}}{R^{2}}\right)dt^{2} + \left(1 + \frac{r^{2}}{R^{2}}\right)^{-1}dr^{2} + r^{2}d\Omega_{2}^{2} \qquad R^{2} \equiv -\frac{3}{\Lambda}$$
(4.23)

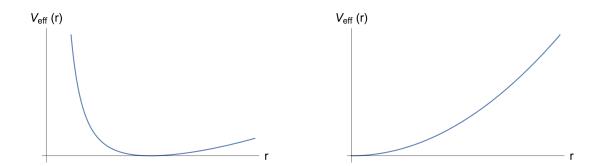


Figure 4.3: Effective potential for anti-de Sitter geodesics, with $\ell \neq 0$ and $\ell = 0$ respectively.

Equivalent coordinates are found setting $r = R \sinh \rho$:

$$ds^{2} = -\cosh^{2}\rho \, dt^{2} + R^{2}d\rho^{2} + R^{2}\sinh^{2}\rho \, d\Omega_{2}^{2}$$
(4.24)

In AdS there's no coordinate singularity in r and indeed now coordinates cover the whole space.

4.2.3.1 Anti-de Sitter geodesics

AdS has the same action as in Eq. 4.15, thus the radial trajectory of a massive particle moving along a geodesic in the $\vartheta = \frac{\pi}{2}$ plane is $\dot{r}^2 + V_{\text{eff}}(r) = E^2$, with:

$$V_{\text{eff}}(r) = \left(1 + \frac{\ell^2}{R^2}\right) \left(1 + \frac{r^2}{R^2}\right) \tag{4.25}$$

Plotting it in Fig. 4.3, the geodesics' behavior is clear: with no angular momentum, anti-de Sitter space acts like a harmonic oscillator, pulling the particle towards the origin and making it oscillate around r = 0. If the particle has non-zero angular momentum, then the potential has a minimum at $r_*^2 = R\ell$, thus particles oscillate around r_* : importantly, particles with finite energy cannot escape to $r \to \infty$, but are constrained by the spacetime within some finite distance from the origin.

This picture of AdS as a harmonic trap which pulls particle to the origin clashes with the fact that AdS is a homogeneous space (roughly, all points are the same). To understand how these two facts are compatible, consider a stationary observer at r=0: this is a geodesic and, from its perspective, other observers (with $\ell=0$) will oscillate around r=0 along geodesics. However, since these observers move along geodesics, in their reference frame they are stationary at r=0, with all other observers oscillating. Thus, while in dS each observer views themselves at the center of the universe, with other observer moving away from them, in AdS each observer views themselves as the center of the universe, with other observer oscillating around them.

Its possible to study geodesics for a massless particle too. This time, the constraint to the action is:

$$-f(r)^{2}\dot{t}^{2} + f(r)^{-2}\dot{r}^{2} + r^{2}(\dot{\vartheta}^{2} + \sin^{2}\vartheta\dot{\varphi}^{2}) = 0$$

This means that the particle follows a null geodesic. Its equation of motion is:

$$\dot{r}^2 + V_{\text{null}}(r) = E^2$$
 $V_{\text{null}}(r) = \frac{\ell^2}{r^2} \left(1 + \frac{r^2}{R^2} \right)$

Its plot in Fig. 4.4 makes it clear that massless particle can escape to $r \to \infty$, as the null potential is asymptotically constant, and it will only experience the usual gravitational redshift. AdS only

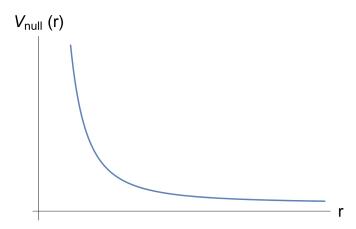


Figure 4.4: Potential experienced by massless particles with $\ell \neq 0$ in AdS.

confines massive particles. To solve for the trajectory, it's easier to work with $r = R \sinh \rho$ and $\ell = 0$, so that:

$$R\dot{\rho} = \pm \frac{E}{\cosh \rho} \quad \Rightarrow \quad R \sinh \rho(\sigma) = E(\sigma - \sigma_0)$$

One sees that $\rho \to \infty$ only in infinite affine time (i.e. $\sigma \to \infty$). However, in coordinate time, recalling Eq. 4.16, $E = \cosh^2 \rho \, \dot{t}$, so:

$$R \tan \frac{t}{R} = E(\sigma - \sigma_0)$$

Hence, as affine time $\sigma \to \infty$, coordinate time $t \to \frac{\pi}{2}R$. This means that light rays escape to infinity in a finite amount of coordinate time: to make sense of dynamics in AdS, one must specify some boundary conditions at infinity to dictate how massless particles and field behave. AdS does not appear to be of any cosmological interest.

4.2.3.2 Anti-de Sitter embeddings

AdS too has a natural embedding in a 5d spacetime: it is a submanifold of $\mathbb{R}^{2,3}$ with metric:

$$ds^{2} = -(dX^{0})^{2} - (dx^{4})^{2} + \sum_{i=1}^{3} (dX^{i})^{2}$$
(4.26)

In particular, anti-de Sitter metric Eq. 4.23 is a metric on the hyperboloid:

$$-(X^{0})^{2} - (X^{4})^{2} + \sum_{i=1}^{3} (X^{i})^{2} = -R^{2}$$
(4.27)

This constraint can be solved via the parametrization:

$$X^{0} = R \cosh \rho \sin \frac{t}{R} \qquad X^{4} = R \cosh \rho \cos \frac{t}{R} \qquad X^{i} = y^{i} R \sinh \rho : \sum_{i=1}^{3} (y^{i})^{2} = 1$$

Given this constraint, $\{y^i\}_{i=1,2,3}$ parametrize \mathbb{S}^2 . The pull-back of the metric Eq. 4.26 on the hyperboloid so parametrized gives anti-de Sitter metric Eq. 4.24.

There's a small subtlety: the embedding hyperboloid has topology $\mathbb{S}^1 \times \mathbb{R}^3$, not \mathbb{R}^4 : this corresponds to a compact time direction, as $t \in [0, 2\pi R)$. However, AdS metric does not have such restriction, with $t \in \mathbb{R}$: this is a universal covering of the hyperboloid.

Another useful parametrization of the hyperboloid is the following:

$$X^4 - X^3 = r$$
 $X^4 + X^3 = \frac{R^2}{r} + \frac{r}{R^2} \eta_{ij} dx^i dx^j$ $X^i = \frac{r}{R} x^i, i = 0, 1, 2$

with $r \in [0, \infty)$. With these coordinates, the metric takes the form:

$$ds^{2} = R^{2} \frac{dr^{2}}{r^{2}} + \frac{r^{2}}{R^{2}} \eta_{ij} dx^{i} dx^{j}$$

These coordinates don't cover the whole AdS, but only one-half of the hyperboloid, restricted to $X^4 - X^3 > 0$: this is known as the *Poincaré patch* of AdS. Moreover, x^0 cannot be further extended beyond $x^0 \in (-\infty, +\infty)$, thus in global coordinates the restriction $t \in [0, 2\pi R)$ remains. In this case, massive particles fall towards r = 0.

Finally, there are two more possible coordinate systems on the Poincaré patch, obtained by setting $z = \frac{R^2}{r}$ and $r = Re^{\rho}$:

$$ds^{2} = \frac{R^{2}}{z^{2}} \left(dz^{2} + \eta_{ij} dx^{i} dx^{j} \right) \qquad ds^{2} = R^{2} d\rho^{2} + e^{2\rho} \eta_{ij} dx^{i} dx^{j}$$

In each case, massive particles fall towards $z = \infty$ or $\rho = -\infty$.

4.3 Symmetries

What makes Minkowski, dS and AdS special solutions to the Einstein equations are their symmetries. The symmetries of Minkowski spacetime are familiar: translations and rotations of spacetime, with the latter splitting between proper rotations and Lorentz boosts. These symmetries are responsible for the existence of energy, momentum and angular momentum on a fixed Minkowski background. It is important to characterize the symmetries of a general metric.

4.3.1 Isometries

Recall Def. 2.3.7 of flow on a manifold: by Eq. 2.10, it is possible to identify a flow with a vector field $K \in \mathfrak{X}(\mathcal{M})$ such that it is tangent to the flow at each point of the manifold: $K^{\mu} = \frac{dx^{\mu}}{dt}$, with $x^{\mu}(t) \equiv x^{\mu}(\sigma_t)$.

Definition 4.3.1. Given a flow associated to $K \in \mathfrak{X}(\mathcal{M})$, it is said to be an *isometry* if:

$$\mathcal{L}_K g = 0 \quad \Leftrightarrow \quad \nabla_{\mu} K_{\nu} + \nabla_{\nu} K_{\mu} = 0 \tag{4.28}$$

Recall Eqq. 4.11-4.12 for the equivalence. This condition means that the metric doesn't change along flow lines: this is called *Killing equation*, and a vector field which satisfied it is a *Killing vector field*. Prop. 2.3.6 can be generalized to hold for the Lie derivative of arbitrary tensor fields, thus Killing vectors too form a Lie algebra: it is the Lie algebra of the isometry group of the manifold.

Example 4.3.1. Given a metric $g_{\mu\nu}(x)$, if it doesn't depend on $x^1 \equiv y$, then $X = \partial_y$ is a Killing vector field, for $(\mathcal{L}_{\partial_y}g)_{\mu\nu} = \partial_y g_{\mu\nu} = 0$. This is the case of ignorable degrees of freedom, as in Eqq. 4.16-4.17.

4.3.1.1 Minkowski spacetime

Consider Minkowski spacetime $\mathbb{R}^{1,3}$ with $g_{\mu\nu} = \eta_{\mu\nu}$. Killing equation is:

$$\partial_{\mu}K_{\nu} + \partial_{\nu}K_{\mu} = 0$$

There two forms of solution. The first one is:

$$K_{\mu} = c_{\mu}$$

for any constant vector c_{μ} . These generate translations. Alternatively:

$$K_{\mu} = \omega_{\mu\nu} x^{\nu} : \omega_{\mu\nu} = -\omega_{\nu\mu}$$

These generate rotations and Lorentz boosts. The emergence of the Lie algebra structure can be elucidated defining the Killing vectors:

$$P_{\mu} := \partial_{\mu} \qquad M_{\mu\nu} := \eta_{\mu\rho} x^{\rho} \partial_{\nu} - \eta_{\nu\rho} x^{\rho} \partial_{\mu}$$

These are 10 Killing vectors: 4 for translations and 6 for rotations and boosts.

Proposition 4.3.1. Killing vectors of Minkowski spacetime obbey:

$$[P_{\mu}, P_{\nu}] = 0 \qquad [M_{\mu\nu}, P_{\sigma}] = -\eta_{\mu\sigma}P_{\nu} + \eta_{\sigma\nu}P_{\mu}$$

$$[M_{\mu\nu}, M_{\rho\sigma}] = \eta_{\mu\sigma} M_{\nu\rho} + \eta_{\nu\rho} M_{\mu\sigma} - \eta_{\mu\rho} M_{\nu\sigma} - \eta_{\nu\sigma} M_{\mu\rho}$$

Proof. By direct calculation.

This is precisely the Lie algebra of the Poincaré group $\mathbb{R}^4 \times SO(1,3)$, i.e. the isometry group of Minkowski spacetime.

4.3.1.2 de Sitter and anti-de Sitter space

The isometries of dS and AdS are best seen from their embeddings. The constraint Eq. 4.21 which defined de Sitter space is invariant under the rotations of $\mathbb{R}^{1,4}$, thus dS inherits the SO(1,4) isometry group. Similarly, the constraint Eq. 4.27 which defines anti-de Sitter space is invariant under rotations of $\mathbb{R}^{2,3}$, so AdS inherits the SO(2,3) isometry group. Note that both these groups are 10-dimensional, as $\mathbb{R}^4 \times \text{SO}(1,3)$.

It's simple to determine the 10 Killing spinors of 5d spacetime:

$$M_{AB} = \eta_{AC} X^C \partial_B - \eta_{BC} X^C \partial_A$$

where X^A , A = 0, 1, 2, 3, 4 are 5d coordinates and η_{AB} is the appropriate Minkowski metric, with signature (-, +, +, +, +) for dS and (-, -, +, +, +) for AdS. In either case, the Lie algebra is that of the appropriate Lorentz group:

$$[M_{AB}, M_{CD}] = \eta_{AD} M_{BC} + \eta_{BC} M_{AD} - \eta_{AC} M_{BD} - \eta_{BD} M_{AC}$$

The embedding hyperbolae are both invariant under these Killing vectors: flows generated by them map points on the hyperbolae to other points on the hyperbolae. Therefore, these Killing vectors are inherited by dS and AdS respectively.

Energy Energy is defined by timelike Killing vectors. In anti-de Sitter space there's no problem finding a timelike Killing vector, as the metric in global coordinates Eq. 4.24 is time-independent, so $K = \partial_t$. However, this changes in de Sitter space.

Considering dS in the static patch with coordinates $r^2 = (X^1)^2 + (X^2)^2 + (X^3)^2$, $X^0 = \sqrt{R^2 - r^2} \sinh \frac{t}{R}$ and $X^4 = \sqrt{R^2 - r^2} \cosh \frac{t}{R}$, the metric Eq. 4.14 is time-independent, thus $K = \partial_t$ is a Killing vector; pushing it forward to the 5d space:

$$\partial_t = \frac{\partial X^A}{\partial t} \partial_A = \frac{1}{R} (X^4 \partial_0 + X^0 \partial_4)$$

On the static patch, this Killing vector is timelike and the energy Eq. 4.16 follows from it. The problem is that the static patch is only one-half of the hyperboloid: when considering the whole AdS, one must account for the case $X^0 = 0, X^4 < 0$, in which the Killing vector points in the past direction, or $X^0 \neq 0, X^4 = 0$, in which it is spacelike. Therefore, the Killing vector can be both timelike (future-directed and past-directed) or spacelike when spanning the whole manifold: for this reason, it's not possible to define a global positive conserved energy on the total de Sitter space. The same conclusion follows by noting that the metric in global coordinates Eq. 4.22 is time-dependent.

4.3.2 Conserved quantities

It is possible to reframe Noether's theorem in the context of Killing vectors.

Theorem 4.3.1 (Noether). Given a massive particle moving on a geodesic $x^{\mu}(\tau)$ in a spacetime with metric g which admits a Killing vector field K, then the Noether charge is conserved along the geodesic:

$$Q := K_{\mu} \frac{dx^{\mu}}{d\tau} \tag{4.29}$$

Proof. First, show that Q is indeed conserved along the geodesic, recalling Eq. 4.28:

$$\frac{dQ}{d\tau} = \partial_{\nu} K_{\mu} \frac{dx^{\nu}}{d\tau} \frac{dx^{\mu}}{d\tau} + K_{\mu} \frac{d^{2}x^{\mu}}{d\tau^{2}} = \partial_{\nu} K_{\mu} \frac{dx^{\nu}}{d\tau} - K_{\mu} \Gamma^{\mu}_{\rho\sigma} \frac{dx^{\rho}}{d\tau} \frac{dx^{\sigma}}{d\tau} = \nabla_{\nu} K_{\mu} \frac{dx^{\nu}}{d\tau} \frac{dx^{\mu}}{d\tau} = 0$$

To show that Q follows from Noether's theorem, consider the action of the massive particle and introduce an infinitesimal transformation $\delta x^{\mu} = K^{\mu}(x)$:

$$\delta \mathcal{S} = \delta \int d\tau \, g_{\mu\nu}(x) \frac{dx^{\mu}}{d\tau} \frac{dx^{\nu}}{d\tau} = \int d\tau \left[\partial_{\rho} g_{\mu\nu} \frac{dx^{\mu}}{d\tau} \frac{dx^{\nu}}{d\tau} \delta x^{\rho} + 2g_{\mu\nu} \frac{dx^{\mu}}{d\tau} \frac{d\delta x^{\nu}}{d\tau} \right]$$

$$= \int d\tau \left[\partial_{\rho} g_{\mu\nu} \frac{dx^{\mu}}{d\tau} \frac{dx^{\nu}}{d\tau} K^{\rho} + 2 \frac{dx^{\mu}}{d\tau} \left(\frac{dK_{\mu}}{d\tau} - K^{\nu} \frac{dg_{\mu\nu}}{d\tau} \right) \right]$$

$$= \int d\tau \left[\partial_{\rho} g_{\mu\nu} K^{\rho} - 2K^{\rho} \partial_{\nu} g_{\mu\rho} + 2 \partial_{\nu} K_{\mu} \right] \frac{dx^{\mu}}{d\tau} \frac{dx^{\nu}}{d\tau} = \int d\tau \, 2\nabla_{\nu} K_{\mu} \frac{dx^{\mu}}{d\tau} \frac{dx^{\nu}}{d\tau}$$

The transformation is a symmetry if $\delta S = 0$, thus, by the symmetry of $\frac{dx^{\mu}}{d\tau} \frac{dx^{\nu}}{d\tau}$, the resulting equation is $\nabla_{(\nu} K_{\mu)} = 0$, i.e. Killing equation.

Example 4.3.2. Energy and angular momentum defined on de Sitter geodesics by Eqq. 4.16-4.17 are Noether charges associated to Killing vectors ∂_t and ∂_{φ} respectively.