# Analysis 3H

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- Adapted from notes of D. Schütz, Durham
- This was part of the Analysis 3H module elective. This is a course on real analysis, touching on metric spaces, tangent spaces, vector fields, manifolds, and differential forms.
- TODO! diagrams, notation (bold vs not bold)

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#### 1.1 Basic notions

The field of real numbers  $\mathbb R$  is a totally ordered field which also satisfies the **completeness** axiom, i.e. a non-empty bounded set  $A \subseteq \mathbb R$  has a **supremum** and/or an **infimum**. The supremum of  $A \subseteq \mathbb R$  is a real number s where  $a \le s$  for all  $a \in A$ . If m is also such that  $a \le m$  for  $a \in A$ , then  $s \le m$ , denoted sup A. The infimum of A is where the inequalities signs are swapped, denoted inf A.

**Lemma 1.1.1** Let  $I_n = [a_n, b_n]$  be a sequence of closed intervals such that  $a_n \le a_{n+1} < b_{n+1} \le b_n$  for all  $n \ge 1$ , then  $\bigcap_{n=1}^{\infty} I_n$  is non-empty.

**Proof** Let  $a = \sup\{a_n\}$ . Since  $a_n \le b_1$  for all n exists by completeness axiom,  $a_n \le b_k$  for any value of n and k, and so  $a \le b_k$ . Hence  $a_k \le a \le b_k$  for all k, and that  $a \in \bigcap_{n=1}^{\infty} I_n$ .

Let M be a set. A function  $d: M \times M \to [0, \infty)$  is called a **metric** on M if

- 1. d(x,y) = 0 iff x = y;
- 2. d(x,y) = d(y,x) for all  $x,y \in M$ ;
- 3.  $d(x,z) \le d(x,y) + d(y,z)$  for all  $x,y,z \in M$ .

The pair (M, d) is then called a **metric space**. It is easy to see any  $N \subseteq M$  is also a metric space using the same d.

**Example** 1. On  $\mathbb{R}$ , d(x,y) = |y - x| gives a metric.

2. On  $\mathbb{R}^2$ ,  $d_1(x,y) = |y_1 - x_1| + |y_2 - x_2|$  is also a metric, but notice that, for example,  $d_1((1,1),(0,0)) = 2$  as opposed to the expected distance of  $\sqrt{2}$ .

The standard (Euclidean) metric in  $\mathbb{R}^2$  is given by

$$d(x,y) = \sqrt{|x_1 - y_1|^2 + |x_2 - y_2|^2}.$$

Let *V* be a real vector space. An **inner product** on *V* is a function  $(\cdot, \cdot) : V \times V \to \mathbb{R}$  that, for all  $x, y \in V$ , satisfies the following:

We will not be distinguishing vectors by bold quantities in this document.

- linearity in the first factor;
- (x,y) = (y,x);
- $(x, x) \ge 0$  and is zero iff x = 0.

**Example** 1. For  $V = \mathbb{R}^n$ , the standard inner product is given by  $(x,y) = x_i y_i$  (where Einstein notation is understood). If A is a symmetric matrix, then  $(x,y) = x^T A y$  is an inner product if all eigenvalues of A are positive.

2. For V = C[a,b],  $(f,g) = \int_a^b f(x)g(x) \, dx$  is an inner product since V is a vector space of continuous functions, and the only function that is everywhere zero and continuous is f(x) = 0 for all  $x \in [a,b]$ .

**Theorem 1.1.2 (Cauchy–Schwartz inequality)** Let V be a real vector space, and  $(\cdot, \cdot)$  an inner product on V. Then

$$|(x,y)| \leq ||x|| \cdot ||y||,$$

where  $\|\cdot\|$  is the standard Euclidean norm of the vector, and there is equality iff  $x = \lambda y$  for some  $\lambda \in \mathbb{R}$ .

**Proof** Note that (x,0) = (x,x-x) = (x,x) - (x,x) = 0, so we may assume that  $y \neq 0$ . Then, with  $\lambda = -(x,y)/\|y\|^2$ ,

$$0 \le (x + \lambda y, x + \lambda y) = ||x||^2 + 2\lambda(x, y) + \lambda^2 ||y||^2$$
$$= ||x||^2 - \frac{(x, y)^2}{||y||^2}.$$

So  $(x,y)^2 \le ||x||^2 ||y||^2$  and the result follows.

**Lemma 1.1.3** Let V be a real vector space with inner product  $(\cdot, \cdot)$ . Then  $d: V \times V \to [0, \infty)$  with d(x, y) = ||x - y|| gives a metric on V.

**Proof** Clearly d(x, x) = 0 and is symmetric, so we just need to check the triangle inequality. By Cauchy–Schwartz,

$$||a+b|| = \sqrt{||a||^2 + 2(a,b) + ||b||^2}$$

$$\leq \sqrt{||a||^2 + 2||a|| ||b|| + ||b||^2}$$

$$\leq ||a|| + ||b||,$$

as required.

Let  $f: M \to N$  be a function metric metric spaces  $(M, d_M)$  and  $(N, d_N)$ . For  $a \in M$ , f is **continuous at** a if, for all  $\epsilon > 0$ , there exists  $\delta > 0$  such that  $d_N(f(a), f(x)) < \epsilon$  for all  $x \in M$  when  $d_M(a, x) < \delta$ .

### 1.2 Sequences and Cauchy sequences

Let M be a metric space. A **sequence**  $(a_n)$  in M consists of elements  $a_n \in M$  for all  $n \in \mathbb{N}$ . Let  $a \in M$ , and  $(a_n)$  **converges to** a if, for all  $\epsilon > 0$ ,  $d(a_n, a) < \epsilon$  for some all  $n \ge n_0$ . We write  $\lim_{n \to \infty} a_n = a$ . The sequence  $(a_n)$  is called **convergent** if there exists  $a \in M$  where  $a_n \to a$ .

**Lemma 1.2.1** Let  $f: M \to N$  be a function between metric spaces and  $a \in M$ . The function f is continuous at  $a \in M$  iff  $f(a_n) \to f(a)$  for  $(a_n) \in M$  with  $a_n \to a$ . (Note that  $f(a_n)$  is a sequence in N.)

**Proof** Assume that f is continuous at  $a \in M$ , and let  $(a_n)$  be a sequence with  $a_n \to a$ . By continuity, for any  $\epsilon > 0$ , there exists  $\delta > 0$  such that, for  $d(a,y) < \delta$ ,  $d(f(a),f(y)) < \epsilon$  for arbitrary  $y \in M$ . Choose  $n_0 \ge 0$  such that  $d(a_n,a) < \delta$  for all  $n \ge n_0$ , then this implies  $d(f(a_n),f(a)) < \epsilon$ , and thus  $f(a_n) \to f(a)$  as required.

On the other hand, assume  $f(a_n) \to f(a)$  for all sequences such that  $a_n \to a$ . Given  $\epsilon > 0$ , assume that instead there is no  $\delta > 0$  such that, for  $d(a,y) < \delta$ ,  $d(f(a),f(y)) < \epsilon$  for arbitrary  $y \in M$ . Then we can find  $a_n \in M$  with  $d(a,a_n) < 1/n$ . However, this means  $d(f(a),f(a_n)) \ge \epsilon$ , which contradicts the assumption that  $f(a_n) \to f(a)$  even though  $a_n \to a$ . So such  $\delta$  exists and we have continuity.

**Lemma 1.2.2** *The limit of a sequence is unique.* 

**Proof** Assume there are two limits a and b for the sequence  $a_n$ . Then  $d(a,b) \le d(a,a_n) + d(a_n,b)$ . As  $n \to \infty$ , the RHS tends to zero so a = b.

A **Cauchy sequence**  $(a_n)$  in the metric space M is a sequence such that, for all  $\epsilon > 0$ , there exists  $n_0 \ge 0$  such that  $d(a_p, a_q) < \epsilon$  for all  $p, q \ge n_0$ .

**Lemma 1.2.3** A convergent sequence is a Cauchy sequence (the converse is not true).

**Proof** Suppose  $a_n \to a$ . Then, for all  $\epsilon > 0$ , there is some  $n_0 \ge 0$  such that  $d(a_n, a) < \epsilon/2$  for  $n \ge n_0$ . Let  $p, q \ge n_0$ , then  $d(a_n, a_q) \le d(a_p, a) + d(a_q, a) < \epsilon$ , so the sequence is Cauchy.

A metric space *M* is **complete** if all Cauchy sequences in *M* converges.

**Theorem 1.2.4** *The real line*  $\mathbb{R}$  *is complete.* 

**Proof** Let  $(a_n)$  be a Cauchy sequence in  $\mathbb{R}$ . Define the sequence of integers  $(n_k)$  where  $n_0 = 1$ , and  $n_{k+1}$  is the smallest integer bigger

than  $n_k$  where  $|a_p-a_q|<2^{-(k+2)}$  for  $p,q\geq n_{k+1}$ . Define the intervals  $I_k=[a_{n_k}-2^{-k},a_{n_k}+2^{-k}]$  and let  $x\in I_{k+1}$ . Now, since  $x\in I_{k+1}$ , this implies that  $|x-a_{n_{k+1}}|<2^{-(k+1)}$ . By definition of the integer sequence,  $|a_{n_k}-a_{n_{k+1}}|<2^{-(k+1)}$ , so then, by triangle inequality,

$$|a_{n_k}-x| \leq |x-a_{n_{k+1}}|+|a_{n_{k+1}}-a_{n_k}| < 2 \cdot 2^{-(k+1)} = 2^{-k},$$

so  $x \in I_k$ . However,  $x \in I_{k+1}$ , so  $I_{k+1} \subset I_k$ . By Lemma 1.1.1,  $\bigcap_{k=1}^{\infty} I_k \neq \emptyset$ , so assume  $a \in \bigcap_{k=1}^{\infty} I_k$ . For  $m \geq n_k$ ,

$$|a - a_m| \le |a - a_{n_k}| + |a_{n_k} - a_m| \le 2^{-k} + 2^{-(k+1)} \to 0$$

as  $m \ge n_k \to \infty$ . Thus  $a_m \to a$  and this arbitrary Cauchy sequence converges in  $\mathbb R$  and thus  $\mathbb R$  is complete.

**Proposition 1.2.5** For  $X \neq \emptyset$ , let  $\mathcal{B}(X)$  be the set of functions  $f: X \to \mathbb{R}$  such that f is bounded. For  $f, g \in \mathcal{B}(X)$ , let  $d(f, g) = \sup_{x \in X} |f(x) - g(x)|$ . Then  $(\mathcal{B}(X), d(f, g))$  defines a complete metric space.

**Proof** d is clearly a metric. For completeness, let  $(f_n)$  be a Cauchy sequence in  $\mathcal{B}(X)$ . For  $x \in X$ ,  $(f_n(x))$  is a Cauchy sequence of real numbers because, by definition of d(f,g),  $|f_q(x) - f_p(x)| \le d(f_p - f_q)$ , and since  $\mathcal{R}$  is complete, the sequence  $(f_n(x))$  converges.

Defining  $f: X \to \mathbb{R}$  such that  $f(x) = \lim_{n \to \infty} f_n(x)$ , we need to show that  $f \in \mathcal{B}(X)$ , and that indeed  $f_n(x) \to f(x)$  regardless of  $x \in X$ . Be definition of a Cauchy sequence, for  $\epsilon > 0$ , there exists  $n_0 \geq 0$  such that  $d(f_p, f_q) < \epsilon/2$  for  $p, q \geq n_0$ . Note also that, for all  $x \in X$ , there exists  $n_1(x) \geq n_0$  such that  $|f_{n_1(x)} - f| < \epsilon/2$ . Then, let  $x \in X$  and  $n \geq n_0$ , we have

$$|f_n(x) - f(x)| \le |f_n - f_{n_1(x)}| + |f_{n_1(x)} - f| < \epsilon.$$

Additionally note,  $|f(x)| \leq |f(x) - f_{n_0(x)}| + |f_{n_0(x)}| \leq \epsilon + c_{f_{n_0}}$  since  $f_{n_0(x)}$  is bounded, so  $f \in \mathcal{B}(X)$ . Further,  $d(f_n - f) = \sup |f_n - f| = \delta < \epsilon$ , so  $f_n$  converges to  $f \in \mathcal{B}(x)$ . Thus every Cauchy sequence converges and thus the space is complete and equipped with a metric.

# Topology of metric spaces

Let (M,d) be a metric space with  $x \in M$  and r > 0. Define the **open** ball around x of radius r to be

$$B(x;r) = \{ y \in M : d(x,y) < r \}.$$

The analogous **closed ball** D(x;r) is defined with the less than or equal to sign. A set  $A \subset M$  is **bounded** if it can be contained in some

D(x;r) for some  $x \in M$ , r > 0. A set  $U \subset M$  is **open** if, for all  $x \in U$ , there exists  $r_x > 0$  such that  $B(x;r_x) \subset U$ . A set  $A \subset M$  is **closed** if  $M \setminus A$  is open.

**Lemma 1.3.1** *Let* (M, d) *be a metric space, then:* 

- 1. M and Ø are open;
- 2.  $\bigcup_i A_i$  is open if all  $A_i \subset M$  are open;
- 3.  $\bigcap_{i=1}^{n} a_{i} = 1$  is open if all  $A_{i} \subset M$  are open and  $n < \infty$ ;
- 4. B(x;r) is open for some r > 0.

**Proof** The first two are obvious. For 3), suppose the open sets  $U_i$  indexed by i are open and  $x \in \bigcap_{i=1}^n U_i$ . Then  $xinU_i$  for all i, so there is some  $B(x;r_i) \subset U_i$ . Taking the minimum of such  $r_i > 0$  means  $B(x;r_i) \subset \bigcap_{i=1}^n U_i$ , and thus the collective finite union is open.

For 4), let 
$$y \in B(x;r)$$
,  $r_y = r - d(x,y) > 0$  and  $z \in B(y;r_y)$ .  
Then  $d(x,z) \le d(x,y) + d(y,z) < d(x,y) + r - d(x,y) = r$ , so  $B(y;r_y) \subseteq B(x;r)$ .

**Corollary 1.3.2** The following may be shown by considering the appropriate complements:

- 1. *M* and ∅ are closed;
- 2.  $\bigcap_i A_i$  is closed if  $A_i \subset M$  for all i;
- 3.  $\bigcup_i A_i$  is closed if  $A_i \subset M$  for all i and  $n < \infty$ ;
- 4. D(x;r) is closed.

**Example** Open intervals are open and closed intervals are closed.

 $(a, \infty)$  is open as it is a union of open bounded intervals.

 $[a, \infty)$  is closed since  $(-\infty, a)$  is open.

 $\mathbb{Z}$  is closed as  $\mathbb{R} \setminus (\bigcup_{n=-\infty}^{\infty} (n, n+1))$  is closed.

Q and [0,1) are neither, while  $\mathbb{R}$  is both.

**Proposition 1.3.3** *Suppose M is a metric space and A*  $\subseteq$  *M. A is closed iff every sequence converges to a*  $\in$  *A.* 

**Proof** Assume A is closed and  $a_n \to a$ . Assume the converse so that  $a \in U = M \setminus A$  which is an open set. Then there is some r > 0 such that  $B(a;r) \in U$ , and since  $a_n \to a$ , there exists  $n_0 \ge 0$  where  $d(a_n,a) < r$  for  $n \ge n_0$ . This implies  $a_n \in B(a;r)$  for all n, but this is a contradiction since  $a_n \in A$ , and thus  $a \in A$ .

Assume  $a_n \to a \in A$ . Let  $x \in M \setminus A$ , r > 0, and assume there is no such  $B(x;r) \subset M \setminus A$ . Thus there is an intersection, i.e.,  $B(x;1/n) \cap A \neq \emptyset$ . This implies that there is some i where  $a_i \in B(x;1/n) \cap A$ . However,  $(a_n)$  is a sequence in A and  $d(a_m,x) < 1/n$  for  $m \ge n+1$ , so  $a_m \to a$ , but this implies x = a which is not possible since  $x \in M \setminus A$ . So  $M \setminus A$  is open which means A is closed.

**Theorem 1.3.4** *Let* M *be a complete metric space and*  $A \subseteq M$  *is closed. Then* A *is complete with the induced metric.* 

**Proof** Let  $(a_n)$  be a Cauchy sequence in A. Since M is complete,  $(a_n)$  converges in M, but A is closed, so  $(a_n)$  converges in A by previous proposition, which implies A is complete.

Let M be a metric space. M is **compact** if every sequence  $(a_n) \in M$  has a convergent subsequence  $(a_{n_k})$ .

**Example** •  $(a_n) = (-1)^n$  is non-convergent but has a convergent sequence.

- M = (0,1) is not compact since  $a_n = 1/n$  and its subsequences do not converge in M.
- $\mathbb{R}$  is not compact as  $a_n$  has no subsequence converging in  $\mathbb{R}$ .
- M=[0,1] is compact. Let  $(a_n)$  be a subsequence in M. Let  $I_1$  be either [0,1/2] or [1/2,1], and let  $(a_{n_k})$  be the subsequences in  $I_1$ . Continuing this we have a sequence of intervals  $I_{m+1} \subset I_m$  with  $I_m$  of length  $2^{-m}$ . Denote the subsequences  $(a_{m_k}^m)$  to be those in  $I_m$ . Taking  $b_m = a_{n_m}^m \in I_m$ , we see that  $b_{m+1} \in I_m$  since  $I_{m+1} \subset I_m$ , so that  $d(b_m,b_q) \leq 2^{-m}$  for  $q \geq m$ . Thus  $(b_m)$  is a Cauchy sequence, which is a subsequence of  $(a_n)$ . Since  $M \subseteq \mathbb{R}$ , M is complete, so  $b_m \to b \in M$ , and thus M is compact.

**Proposition 1.3.5** By extension, closed n-gons in  $\mathbb{R}^n$  are compact.

**Proposition 1.3.6** *Let*  $f: M \to N$  *be a continuous map between metric spaces. If* M *is compact, then*  $f(M) \subset N$  *is compact.* 

**Proof** Let  $(a_n)$  be a sequence in f(M). Then  $a_n = f(b_n)$  for some  $b_n \in M$ . The sequence  $(b_{n_k})$  converges in M since M is compact, thus

$$\lim_{k \to \infty} a_{n_k} = \lim_{k \to \infty} f(b_{n_k}) = f\left(\lim_{k \to \infty} b_{n_k}\right) = f(b)$$

since f is continuous. So  $(a_{n_k})$  is convergent, thus f(M) is compact.  $\square$ 

**Proposition 1.3.7** A closed subset of a compact space is a compact set.

**Proof** Let  $(a_n)$  be a sequence in  $A \subset M$  where M is compact. Since  $(a_n) \in M$ ,  $(a_{n_k})$  is convergent, but A closed so  $(a_{n_k}) \to a \in A$ , thus A is compact.  $\square$ 

#### 1.3.1 Heine–Borel theorem

**Theorem 1.3.8** A subset  $A \subseteq \mathbb{R}^n$  is compact iff A is closed and bounded.

**Proof** Suppose A is compact, so clearly A is closed. If A is unbounded, then there exists  $(a_n) \in A$  where  $d(a_n, 0) \ge n$ , so  $(a_{n_k})$  does not converge in  $\mathbb{R}^n$ . However A is compact, which is a contradiction, so A is bounded.

Suppose *A* is bounded, then  $A \subseteq [a, b]^n$ . If *A* is closed, then it is a closed subset of a compact set, so *A* is compact by previous proposition.

For example, if  $f: M \to N$  with f is a scalar continuous function, then  $f(M) \subset \mathbb{R}$  is closed and bounded since M is compact, and thus f(M) compact implies f(M) is closed and bounded.

## 1.4 Banach and Hilbert spaces

Let *V* be a real vector space. The **norm** on *V* is a function  $\|\cdot\|: V \to [0, \infty)$  where:

- 1. ||x|| = 0 iff x = 0;
- 2.  $\|\lambda x\| = |\lambda| \cdot \|x\|$  for all  $x \in V$  and  $\lambda \in \mathbb{R}$ ;
- 3.  $||x + y|| \le ||x|| + ||y||$ .

The pair  $(V, \|\cdot\|)$  gives a **normed vector space**.

**Lemma 1.4.1** Let V be a normed vector space, then d(x,y) = ||x - y|| defines a metric on V.

**Proof** Two of the properties follow from definition. To show the reflexive property, note that

$$d(y,x) = ||y-x|| = ||(-1)(x-y)|| = ||x-y|| = d(x,y).$$

**Example** 1. It may be shown that the metrics

$$\sum_{i} |x_i|, \qquad \sum_{i} \sqrt{|x_i|^2}, \qquad \max\{|x_i| \in \mathbb{R}\}$$

define norms on  $\mathbb{R}^n$  (the  $\ell^1$ ,  $\ell^2$  and  $\ell^\infty$  norms).

**2**. The **supremum norm** on B(X) is defined by

$$||f||_{\infty} = \sup\{|f(x)| \in \mathbb{R} ; x \in X\}.$$

3. For X a metric space,  $C_b(X) = \{f : x \to \mathbb{R} : f \text{ continuous and bounded} \}$  is also a normed vector space with the supremum norm.

If  $C(X) = \{f : x \to \mathbb{R} : f \text{ continuous}\}$  then f does not have a supremum, however, we have the following:

**Proposition 1.4.2** *If* X *is compact, then*  $C(X) = C_b(X)$ *, so* C(X) *is a normed vector space.* 

**Proof**  $C_b(X) \subseteq C(X)$  regardless of X. For the converse, assume  $f \in C(X)$ , so that f(X) is compact. This implies f(X) is bounded and closed by the Heine–Borel theorem, so  $C(X) \subseteq C_b(X)$ .

Let  $(V, \|\cdot\|_V)$  and  $(W, \|\cdot\|_W)$  be two normed vector spaces. A function  $f: V \to W$  is continuous at  $x \in V$  if, for all  $\epsilon > 0$ , there exists  $\delta > 0$  such that  $\|x - y\|_V < \delta$  implies that  $\|f(x) - f(y)\|_W < \epsilon$ .

Let V be a normed vector space. V is a **Banach space** if V with the metric induced by the norm is complete.

**Theorem 1.4.3** Let X be a metric space, then  $C_b(X)$  with the supremum norm is a Banach space.

**Proof** Since  $C_b(X) \subseteq B(X)$ , if  $C_b$  is closed, then  $C_b$  is complete since B(X) is complete. To show this, let  $(f_n) \in C_b(X)$ , and let  $f_n \to f \in B(X)$ . The convergene of  $f_n$  implies that there exists  $n_0 \ge 0$  such that  $\|f_n - f\| < \epsilon/3$  for any  $\epsilon > 0$  with  $n \ge n_0$ . Also,  $\|f_{n_0}(y) - f(y)\| < \epsilon/3$  for all  $y \in X$ . The functions are continuous, so there exists  $\delta > 0$  where, if  $d(x,y) < \delta$ ,  $\|f_{n_0}(x) - f_{n_0}(y)\| < \epsilon/3$  for  $x \in X$ . Thus, for  $d(x,y) < \delta$ ,

$$|f(x)-f(y)| \le |f(x)-f_{n_0}(x)|+|f_{n_0}(x)-f_{n_0}(y)|+|f_{n_0}(y)-f(y)| < \epsilon$$

so f is continuous, and  $C_b(X)$  is closed and thus complete.

**Corollary 1.4.4** For a < b, C[a,b] with the supremum norm is a Banach space.

Note that C[a, b] is not a complete space with, for example, the  $L_2$  norm

$$||f||_2 = \sqrt{\int_a^b (f(x))^2 dx}.$$

For example, with  $f_n = x^n$ ,  $f_n \to 0$  but clearly  $f_n(1) = 1$  for all n. The underlying reason is the sequence is not a Cauchy sequence with respect to the norm.

Convergence with respect to the supremum norm is called **uniform convergence** (cf. Complex Analysis 2H).

Let  $(V, \|\cdot\|)$  be a Banach space. If there is an inner product from V which induces this norm, then V is called a **Hilbert space**.

**Theorem 1.4.5** Let (M,d) be a metric space. Then there exists  $(\overline{M},\overline{d})$  where  $\overline{M}$  is complete, and there is an embedding  $\iota: M \to \overline{M}$  with  $d(x,y) = d(\iota(x),\iota(y))$  for all  $x,y \in M$ . Also, for all  $\overline{x} \in \overline{M}$ , there is a sequence  $(x_n) \in M$  with  $x_n \to \overline{x}$  as  $n \to \infty$ .

Here,  $\overline{M}$  is called the **completion** of M, and it is unique up to some isomorphism.

**Example** The completion of  $\mathbb Q$  is  $\mathbb R$  with respect to the Euclidean metric.

The completeness of C[a, b] with respect to the inner product metric is denoted  $L^2[a, b]$ ..

### 1.4.1 The contraction mapping theorem

**Theorem 1.4.6** Let (M,d) be a complete metric space,  $0 \le \lambda \le 1$  and a  $f: M \to M$  with  $d(f(x), f(y)) \le \lambda d(x, y)$  for all  $x, y \in M$ . Then f has one unique fixed point where  $f(x_0) = x_0$ .

**Proof** Note that f is a contraction, and continuity is automatically satisfied from the condition that  $d(f(x), f(y)) \le \lambda d(x, y)$ .

Let  $x \in M$ , and  $a_n = f^n(x)$ . So we have

$$d(x, a_n) \leq d(x, f(x)) + d(f(x), f^2(x)) + \dots d(f^{n-1}(x), f^n(x))$$

$$= \sum_{i=0}^{n-1} d(f^i(x), f^{i+1}(x))$$

$$\leq \sum_{i=0}^{n-1} \lambda d(x, f(x))$$

$$= d(x, f(x)) \frac{1 - \lambda^n}{1 - \lambda}$$

$$\leq \frac{d(x, f(x))}{1 - \lambda},$$

by Cauchy–Schwartz and the arithmetic progression with  $0 \le \lambda < 1$ . Now,

$$d(a_n, a_m) = d(f^n(x), f^m(x)) \le \lambda^m d(f^{n-m}, x) \le \lambda^m \frac{d(x, f(x))}{1 - \lambda}$$

assuming n > m. For  $n, m \ge n_0$ , we have

$$d(a_n, a_m) \leq \lambda^{n_0} \frac{d(x, f(x))}{1 - \lambda}.$$

Clearly  $(a_n)$  is a Cauchy sequence, and thus we have completeness and  $a_n \rightarrow a \in M$ . Now,

$$f(a) = f\left(\lim_{n \to \infty} a_n\right) = \lim_{n \to \infty} f(a_n) = \lim_{n \to \infty} a_{n+1} = a,$$

Note that elements of  $L^2$  are not exactly functions, but rather *equivalence classes* (cf.  $11 \equiv 1 \mod 10$ )

Or, if you throw a map of the world on the floor, there is exactly one point on the map that exactly corresponds to one point on the floor. so there is some  $a \in M$  that is a fixed point.

To show uniqueness, suppose b is another fixed point. Then

$$d(a,b) = d(f(a), f(b)) \le \lambda d(a,b),$$

and for  $\lambda \neq 0$ , d(a, b) = 0, so a = b.

# 1.5 A norm for matrix spaces

We want a norm reflecting the fact that matrices can be identified with linear maps. Let  $A = (A_{ij}) \in M_{n,k}(\mathbb{R})$ . We define

$$||A|| = \sup\{||Ax||_2 : x \in \mathbb{R}^k, ||x||_2 \le 1\},$$
 (1.1)

where  $\|\cdot\|$  is the Euclidean norm. Here,  $Ax \in \mathbb{R}^n$ , and  $x \mapsto \|Ax\|_2$  is clearly a continuous map. By the Heine–Borel theorem,  $\{\|Ax\|_2 : \|x\|_2 \le 1\}$  is bounded and closed, so the supremum exists, and there is x with  $\|x\|_2 \le 1$  such that  $\|A\| = \|Ax\|_2$  exists.

### Lemma 1.5.1 We have

- $||Ax||_2 \le ||A|| ||x||_2$  for all A and x
- $||AB|| \le ||A|| ||B||$
- $||A||_{\infty} \le ||A|| \le k\sqrt{n} ||A||_{\infty}$ ,

where 
$$||A||_{\infty} = \max\{|A_{ij}| : A \in M_{n,k}(\mathbb{R})\}.$$

Let  $U \subset \mathbb{R}^n$  be open. A **vector field** or **autonomous differential equation** is a continuous map  $v: U \to \mathbb{R}^n$  with no explicit time dependence. Here, U is called the **phase space** of v.

For  $x \in U$ ,  $\tau \in \mathbb{R}$ , a continuous differential curve  $\alpha : (a,b) \to U$  is an **integral curve** of v at  $(x,\tau)$  if  $\tau \in (a,b)$ ,  $\alpha(t) = x$  and  $\alpha'(t) = v(\alpha(t))$ . Note the integral curves have tangent vectors which agree with v at a given point.

More generally, for  $U \in \mathbb{R}^n$ ,  $I \subset \mathbb{R}$ , a **differential equation** is a continuous map  $V: U \times I \to \mathbb{R}^n$ . A **solution** of V at  $x \in U$  and  $\tau \in I$  is a continuously differential curve  $\alpha: I \to U$  with  $\alpha'(t) = V(\alpha(t), t)$  and  $\alpha(t) = x$ .

# 2.1 Picard–Lindelöf theorem

This is an existence and uniqueness theorem for differential equations.

**Theorem 2.1.1** Let  $U \subset \mathbb{R}^n$ ,  $I \subset \mathbb{R}$  be open and  $V: U \times I \to \mathbb{R}^n$  be a differential equation where, for all  $x_1, x_2 \in U$ ,  $t \in I$ , there exists  $L \geq 0$  such that

$$||v(x_1,t)-v(x_2,t)|| \le L||x_1-x_2||_2.$$

Given  $(u, \tau) \in U \times I$ , there exists a, b > 0 with

$$U_1 = \{x \in U : ||x - u|| < a\}, \quad I_1 = \{t \in I : |t - \tau| < b\}$$

such that the differential equation v has an unique solution for all  $x \in U_1$  and  $\tau \in I_1$ . Furthermore, the resulting  $\alpha : U_1 \times I_1 \to U$  given by  $\alpha(x,t) = \alpha_x(t)$  is continuous.

**Proof** This one is quite long! The key idea is to construct a contraction mapping *A* and make use of the fixed point theorem to demonstrate existence and uniqueness. We are going to split this up into little bits.

• We first construct an integral curve  $\alpha$  with  $\partial \alpha / \partial t(x,t) = v(\alpha,t)$ ,

Compare this with the **Lipschitz condition** where  $||v(x_1) - v(x_2)|| \le L||x_1 - x_2||$ , where L is the **Lipschitz constant**.

 $\alpha(x,\tau) = x$ . By integrating,

$$\alpha(x,t) = x + \int_{\tau}^{t} v(\alpha(x,s),s) \, \mathrm{d}s.$$

Define some operator A such that

$$A\beta(x,t) = x + \int_{\tau}^{t} v(\beta(x,s),s) \, \mathrm{d}s,$$

then we note that  $A\alpha = \alpha$ , and  $\alpha$  is a fixed point of the operator A. We aim to show that A is a contraction in a space satisfying the relevant properties.

• Let  $a_1, b_1 > 0$  be such that

$$D_1 = D(u; 2a_1) \subset U$$
,  $D_2 = D(\tau; b_1) \subset I$ .

By the Heine–Borel theorem,  $D_1 \times D_2 \subset \mathbb{R}^{n+1}$  is compact, and so there exists some  $K \ge 0$  such that, with respect to the Euclidean norm, ||v(x,t)|| < K for all  $(x,t) \in D_1 \times D_2$ .

Let a, b > 0 be such that

$$0 < a < a_1, \qquad b < \min\left\{b_1, \frac{a}{K}, \frac{1}{L}\right\}.$$

Recall that  $U_1 = B(u; a)$  and  $I_1 = B(\tau; b)$ , so let

$$M = \{ \beta : U_1 \times I_1 \to D \subset \mathbb{R}^n \}$$

where  $\beta$  is continuous and  $\beta(x,\tau) = x$  for all  $x \in U_1$ . This implies that

$$M \subseteq (C_h(U_1 \times I_1))^n$$

and since  $(C_b(U_1 \times I_1))^n$  is a Banach space with the supremum norm, if *M* is closed, then *M* is complete.

- Suppose  $(\beta_n) \in M$  where  $\beta_n \to \beta$ . For  $(x,t) \in U_1 \times I_1$ ,  $||\beta(x,t) \beta|| = 0$  $|\beta_n(x,t)|| \leq ||\beta - \beta_n||$  so  $|\beta_n| \rightarrow \beta$ , but since  $D_1$  is closed,  $\beta \in D$ and obviously  $\beta_n((x,\tau) \to \beta(x,\tau) = x$ , so M is closed and so is complete.
- If we now consider  $A\beta$ , then we have  $A\beta(x,\tau) = x$  and that

$$||A\beta(x,t) - u|| \le ||A((x,t) - x|| - ||x - u||$$

$$\le \int_{\tau}^{t} ||v(\beta(x,s),s)|| \, ds + a$$

$$\le K|t - \tau| + a$$

$$\le Kb + a$$

$$< 2a < 2a_1$$

by Cauchy–Schwartz, definition of  $U_1$ , second bullet point, definition of  $I_1$ , and definition of b and a respectively. By definition of  $D_1$ , we have  $A\beta(x,t) \in D_1$ .

Recall D denotes closed balls, while B denote open balls.

• Note then we have

$$||A\beta(x,t) - A\beta(y,t')| \le ||x - y|| + \left\| \int_{\tau}^{t} v(\beta(x,s),s) - v(\beta(y,s),s) \, ds \right\|$$

$$+ \left\| \int_{t}^{t'} v(\beta(y,s),s) \, ds \right\|$$

$$\le ||x - y|| + L \int_{\tau}^{t} ||v(\beta(x,s),s) - v(\beta(y,s),s)|| \, ds$$

$$+ K|t - t'|$$

$$\le ||x - y|| + L \sup_{s \in [\tau,t]} ||\beta(x,s) - \beta(y,s)|| + K|t - t'|,$$

by the Lipschitz conditions. All terms can be made arbitrarily small since x can be made close to y, t can be made close to t', and since  $[\tau,t]$  is compact,  $\|\beta(x,s)-\beta(y,s)\|$  can be made arbitrarily small. So now  $A\beta \in D_1$  is continuous, and therefore  $A\beta \in M$ , and  $A:M\to M$  is a self mapping.

• Since *A* is a self-mapping, for  $\beta_{1,2} \in M$ , we have

$$||A\beta_{1} - A\beta_{2}|| \leq \int_{\tau}^{t} ||v(\beta_{1}(x,s),s) - v(\beta_{2}(x,s),s)|| \, ds$$

$$\leq L \int_{\tau}^{t} ||\beta_{1} - \beta_{2}|| \, ds$$

$$= L|t - \tau|||\beta_{1} - \beta_{2}||$$

$$\leq (Lb)||\beta_{1} - \beta_{2}||$$

by definition of  $I_1$ . Note that Lb < 1 by the definition of b, and therefore A is a contraction.

b<1/L.

Since A is a contraction and M is complete, by contraction mapping there is one unique point in M that is fixed under A. Clearly this is  $\alpha$  by definition of  $\beta$  (see first bullet point), and hence  $\alpha$  is the unique solution to the ODE satisfying the stated conditions.

Note that it doesn't matter if  $\alpha: I_1 \to U$ , since we can redefine M and A as  $M_x = \{\beta: I_1 \to D\}$  with  $\beta(t) = x$ , and  $A_x: M_x \to M_x$ . There will be an unique solution for fixed  $x \in U_1$ , where the generation solution gives this solution.

# Differentiation in $\mathbb{R}^n$

2.2

Let  $U \subset \mathbb{R}^n$  be open. Recall that  $f: U \to \mathbb{R}^n$  is differentiable at  $x \in U$  with derivative

$$Df(x) = \left(\frac{\partial f_i}{\partial x_j}\right) \in M_{p,n}(\mathbb{R})$$
 (2.1)

$$f(x+h) = f(x) + Df(x) \cdot h + R(h), \qquad \lim_{\|h\| \to 0} \frac{R(h)}{\|h\|} = 0.$$

If f is differential for all  $x \in U$ , then  $Df : U \to M_{p,n}(\mathbb{R}) = \mathbb{R}^{pn}$ . If  $D^i f$  is continuous then f is said to be of i-class, with  $f \in C^i(U)$ .

#### 2.2.1 Mean value theorem

**Theorem 2.2.1** Let  $U \subset \mathbb{R}^n$  be open,  $x \in U$ ,  $h \in \mathbb{R}^n$  where  $x + th \in U$  for all  $t \in [0,1]$  and  $f \in C^1 : U \to \mathbb{R}^p$ , then

$$f(x+h) - f(x) = \int_0^1 Df(x+th) \cdot h \, dt.$$

**Proof** Let  $f_i: U \to \mathbb{R}$  with  $g_i(t) = f_i(x+th)$ , so that  $g: [0,1] \to \mathbb{R}$ . Then we have  $g_i'(t) = Df_i(x+th) \cdot h$ . By the fundamental theorem of calculus,

$$g_i(1) - g_i(0) = \int_0^1 Df_i(x+th) \cdot h \, dt$$
  
=  $f_i(x+th) - f_i(x)$ .

Since this is true per component, we have the result in higher dimensions.

**Corollary 2.2.2** Let  $U \subset \mathbb{R}^n$  be open and  $convex^1$ , and also that  $f \in C^1$ :  $U \to \mathbb{R}^n$ . Assume that there exists some  $C = \sup\{\|Df(x)\| \in \mathbb{R} : x \in U\}$ , then  $\|f(y) - f(x)\| \le C\|y - x\|$ .

<sup>1</sup> So for all  $x, y \in U$ ,  $xt + (1 - t)y \in U$  for  $t \in [0, 1]$ .

**Proof** By the mean value theorem, we have

$$||f(x+h) - f(x)|| \le \int_0^1 ||Df(x+h \cdot h)|| \, dt$$

$$\le \int_0^1 ||Df(x+h)|| \cdot ||h|| \, dt$$

$$\le \int_0^1 C \cdot ||h|| \, dt = C \cdot ||h||.$$

Since h is arbitrary (up to us assuming convexity), letting h = y - x leads the result.

Note that for the above corollary, U can always be reduced so that C exists locally. For the Picard–Lindelöf theorem, we get  $v \in C^1: U \times I \to \mathbb{R}$  implies the Lipschitz condition is satisfied locally.

#### 2.2.2 *Matrices*

Let  $U \subset \mathbb{R}^n$  and  $V \subset \mathbb{R}^p$  be open. A  $C^1$ -function  $f: U \to V$  is a **diffeomorphism** if there exists  $f^{-1}: V \to U$  where  $f \circ f^{-1} = f^{-1} \circ f = \mathrm{id}$  (the identity map), and  $f^{-1}$  is differential for all  $x \in V$ .

**Example**  $f = x^3$  has  $f^{-1} = x^{1/3}$ , but since  $f^{-1}$  is not differentiable at x = 0,  $x^3$  is not a diffeomorphism on  $\mathbb{R}$ .

By the chain rule, note that

$$D(f^{-1} \circ f) = (Df^{-1}(f))Df = I_n, \quad D(f \circ f^{-1}) = (Df(f^{-1}))Df^{-1} = I_p.$$

If y = f(x) then  $Df^{-1}(y) = (Df(x))^{-1}$ , then inverse matrix of Df(x), so Df(x) is invertible and p = n if f is a diffeomorphism.

**Lemma 2.2.3** 1.  $GL_n(\mathbb{R})$  is an open set.

- 2.  $A \in M_{n,n}(\mathbb{R})$  with  $||A|| \leq 1$  implies that  $I A \in GL_n(\mathbb{R})$ .
- 3.  $inv : GL_n(\mathbb{R}) \to GL_n(\mathbb{R})$  with  $A \mapsto A^{-1}$  is a smooth diffeomorphism.

**Proof** Recall that the determinant is defined as

$$\det A = |A| = \sum_{\mathbf{sig} \in S_n} \operatorname{sign}(\sigma) \prod_{i=1}^n a_{i,\sigma(i)}.$$

This is a polynomial in components of *A*, so it is a smooth function.

- 1.  $GL_n(\mathbb{R}) = \det^{-1}(\mathbb{R} \{0\})$  so  $A \in GL_n(\mathbb{R})$  implies that  $|A| \neq 0$ , which implies  $|B| \neq 0$  for B close to A, and thus  $GL_n(\mathbb{R})$  is open.
- 2. If  $||A|| \le 1$ , define  $B_n = \prod_{i=0}^n A^i$  where  $A^0 = I$ .  $\{B_n\}$  is a Cauchy sequence since

$$||B_n - B_m|| \le \sum_{k=\min\{m,n\}+1}^{\max\{m,n\}} ||A||^k \le \frac{||A||}{1 - ||A||} \to 0$$

for sufficiently large m, n with  $||A|| \le 1$ . So there exists  $B = \lim_{n\to\infty} B_n$ , and thus

$$(I-A)B = (I-A)\lim_{n\to\infty} B_n.$$

 $B_n$  continuous implies that

$$(I-A)B = \lim_{n \to \infty} (I-A)B_n = \lim_{n \to \infty} I - A^{n+1} = I$$

since 
$$||A|| \le 1$$
, so  $B^{-1} = I - A \in GL_n(\mathbb{R})$ .

3. By Cramer's rule, for  $A = (a_{ij})$ ,  $A^{-1} = (b_{ij})$  with  $b_{ij} = \det A_{ij}/\det A$ , where  $A_{ij}$  is the matrix obtained by replacing the  $i^{th}$  column with the standard  $j^{th}$  basis vector. So  $(b_{ij})$  depends smoothly on  $(a_{ij})$  since det is a smooth map, and so inv is smooth. Note additionally that inv  $\circ$  inv = id, so it is a bijection and hence a diffeomorphism.

This is the general linear group with real entries.

 $S_n$  here is the group of symmetric permutations, and  $\operatorname{sig}(\sigma)$  is the signature of the permutation  $\sigma$  (+1 if even and -1 if odd).

#### *Inverse function theorem* 2.2.3

Let  $U \subset \mathbb{R}^n$  be open and  $f \in C^k : U \to \mathbb{R}^n$ . f is locally invertible at  $x \in U$  if there exists  $U_1 \subset U$  such that for  $x \in U_1$ ,  $V_1 \subset \mathbb{R}^n$  where  $f(x) \in V_1$  is open and  $f: U_1 \to V_1$  is a diffeomorphism.

**Theorem 2.2.4** Let  $U \subset \mathbb{R}^n$  be open and  $f \in C^k : U \to \mathbb{R}^n$ ,  $u \in U$ . f is locally invertible iff Df(u) is invertible. Here the local inverse is of class  $C^k$ .

**Proof** This one is quite long!

- If f is locally invertible at u, then it is a diffeomorphism, so clearly Df(u) is invertible. However, this is for an isolated point, and we need to show that is also true on the appropriate neighbourhood.
- Assume that u = 0 = f(u), i.e. a fixed point, and Df(0) = I. Define, for  $y \in \mathbb{R}^n$ ,

$$g_y(x) = y + x - f(x) \quad \Rightarrow \quad y - f(x) = g_y - x.$$

Note that  $Dg_y(x) = I - Df(x)$  and does not depend on y. Also that  $Dg_{y}(0) = I - I = 0$ .

By continuity, we have  $||Dg_y(x)|| = ||Dg_0(x)|| \le 1/2$  for some xnear 0. This implies that

$$\|g_y(x_1 - g_y(x_2))\| \le \frac{1}{2} \|x_1 - x_2\|$$

for  $x_{1,2} \in D(0;r)$ . Taking  $x_2 = 0$ , we also get

$$\|g_y(x) - y\| \le \frac{1}{2} \|x\|,$$

so we have

$$||g_y(x)|| \le \frac{1}{2}||x|| + ||y||$$

for  $y \in D(0;r/2)$  and  $x \in D(0;r)$ , and thus  $||g_y(x)|| \le r$ . Hence we have  $g_y(x): D(0;r) \to D(0;r)$ , and  $g_y(x)$  is by construction a contraction since  $||g_y(x_1 - g_y(x_2))|| \le (1/2)||x_1 - x_2||$ .

• By contraction mapping theorem, for all  $y \in D(0; r/2)$ , there exists a unique  $x \in D(0;r)$  with y = f(x), so there exists an inverse function defined on D(0; r/2).

Define

$$U_1 = \{x \in U : ||x|| < r, ||f(x)|| < r/2\}, \quad V_1 = f(U_1) = B(0; r/2).$$

By definition, both the domain and image are open sets.  $f: U_1 \rightarrow$  $V_1$  is a restricted bijection since it is a bijection on  $D(0;r/2) \supset$ 

B(0; r/2). Given  $x_{1,2} \in D(0; r)$ , we have

$$||x_1 - x_2|| = ||g_0(x_1) + f(x_1) - g_0(x_2) + f(x_2)||$$

$$\leq ||g_0 - g_0(x_2)|| + ||f(x_1) - f(x_2)||$$

$$\leq \frac{1}{2}||x_1 - x_2|| + ||f(x_1) - f(x_2)||,$$

so that  $||x_1 - x_2|| \le 2||f(x_1) - f(x_2)||$ . For  $x_2 = 0$ , we have  $||x_1|| \le 2||f(x_1)||$ . Since  $||f(x_1)|| < r/2$  by construction, we have  $||x_1|| < r$ , so indeed  $V_1 = B(0; r/2)$ .

For  $f^{-1} = \phi$ ,  $||x_1 - x_2|| \le 2||f(x_1) - f(x_2)||$  implies that  $||\phi(y_1) - \phi(y_2)|| \le 2||x_1 - x_2||$ , so that  $f^{-1}$  is Lipschitz continuous.

- Note that Df(x) is invertible for all  $x \in D(0;r)$ , since we have  $g_0(x) x = f(x)$ , so that  $Df(x) = I Dg_0(x)$ , but  $\|Dg_0(x)\| \le 1/2$  from point 2 above, so Df(x) is invertible for all  $x \in D(0;r)$ , and in particular for  $x \in B(0;r) \subset D(0;r)$ .
- Recall that if f id differentiable, then  $f(x_1) f(x_2) = Df(x_1)(x_1 x_2) + R(x_1 x_2)$  with  $R(h)/\|h\| \to 0$  as  $\|h\| \to 0$ . Let  $y_i = f(x_i)$ . For i = 1, 2,

$$y_1 - y_2 = Df(x_1) (\phi(y_1) - \phi(y_2)) + R (\phi(y_1) - \phi(y_2)),$$

so that

$$(Df(\phi(y_1)))^{-1} (y_1 - y_2) = (\phi(y_1) - \phi(y_2)) + (Df(\phi(y_1)))^{-1} R (\phi(y_1) - \phi(y_2)).$$

We want to show that the remainder term tends to zero, which will show that  $\phi = f^{-1}$  is differentiable. For that, note we have, by Cauchy–Schwartz and point 3 above,

$$\frac{\| \left( Df(\phi(y_1)) \right)^{-1} R \left( \phi(y_1) - \phi(y_2) \right) \|}{\| y_1 - y_2 \|} \le \frac{\| \left( Df(\phi(y_1)) \right)^{-1} \| \cdot \| R \left( \phi(y_1) - \phi(y_2) \right) \|}{(1/2) \| \left( \phi(y_1) - \phi(y_2) \right) \|}.$$

 $(Df(\phi(y_1)))^{-1}$  is bounded since f is differentiable. Further more, f differentiable means  $\|R\left(\phi(y_1)-\phi(y_2)\right)\|/\|\left(\phi(y_1)-\phi(y_2)\right)\|\to 0$  as  $\|\left(\phi(y_1)-\phi(y_2)\right)\|\to 0$ . Thus the desired remainder goes to zero since  $y_1-y_2\to 0$  implies  $\phi(y_1)-\phi(y_2)\to 0$ , and  $\phi=f^{-1}$  is differentiable.

• The derivative  $D\phi(y) = (Df(\phi(y)))^{-1} = \text{inv} \circ Df \circ \phi)y$ , so by construction,  $D\phi = Df^{-1}$  is continuous. By chain rule, if  $f \in C^k$ ,  $D^{k-1}\phi$  is continuous, and thus  $\phi = f^{-1} \in C^k$ .

#### *Implicit function theorem* 2.2.4

**Theorem 2.2.5** Let  $U \subset \mathbb{R}^n$ ,  $V \subset \mathbb{R}^m$  be open and  $f: U \times V \to \mathbb{R}^m$ be a  $C^k$ -function, with  $k \geq 1$ . Let  $(u,v) \in U \times V$  such that the matrix  $[\partial f_i/\partial x_i](u,v)$  is invertible with c=f(u,v). Then there is a  $C^k$ -function  $\eta: U_1 \to V_1$  with  $u \in U_1 \subset U$ ,  $v \in V_1 \subset V$  where  $\eta(u) = v$  and  $f(x,\eta(x)) = c$  for all  $x \in N((u,v);r)$ . Further more, if f(x,y) = c for  $(x,y) \in U_1 \times V_1$ , then we have  $y = \eta(x)$  in the respective sets.

**Proof** Define  $\phi: U \times V \to \mathbb{R}^n \times \mathbb{R}^m$  where  $(x,y) \mapsto (x,f(x,y))$ . We have

$$D\phi(u,v) = \begin{pmatrix} I & 0 \\ \partial f_i/\partial x_j(u,v) & \partial f_i/\partial x_j(u,v) \end{pmatrix},$$

so  $\det D\phi(u,v) \neq 0$ , and so by the inverse function theorem,  $\phi$  is locally a diffeomorphism.

Since  $\phi(x,y) = (x, f(x,y))$ , we have  $\phi^{-1}(A,b) = (A,g(A,b))$ . Setting  $\eta(x) = g(x,c)$ , then defining  $\hat{p}_2$  as the projection operator for the second argument, we have

$$f(x,\eta(x)) = f(x,g(x,c))$$

$$= \hat{p}_2 \phi(x,g(x,c))$$

$$= \hat{p}_2 \phi \phi^{-1}(x,c)$$

$$= \hat{p}_2(x,c) = c.$$

So we have f(x,y) = c iff  $y = \eta(x)$  for  $(x,y) \in W$  where  $\phi$  is a diffeomorphism. This is achieved by choosing  $u \in U_1 \subset U$ ,  $v \in V_1 \subset$ *V* so that  $U_1 \times V_1 \subset W$ , with  $\eta(U_1) = V_1$ .

The implicit function theorem gives a criterion of when we can solve f(x,y) = c unique for y. In fact, if the linear equation [Df(u,v)](x,y) =0 is uniquely solvable, then f(x,y) = c is uniquely solvable for y.

#### 2.2.5 Manifolds

Let  $M \subset \mathbb{R}^n$ ,  $k \geq 0$ ,  $\ell \geq 1$ . M is a  $C^{\ell}$  k-dimensional manifold if, for all  $p \in M$ , we also have  $p \in U \subset \mathbb{R}^n$  where there exists a  $C^{\ell}$ -diffeomorphism  $h: U \to U' \subset \mathbb{R}^k \times \mathbb{R}^{n-k}$  with  $h(U \cap M) =$  $U' \cap (\mathbb{R}^k \times \{0\})$ . Informally, a manifold is a structure where every point of M has a neighbourhood that resembles  $\mathbb{R}^k$ . h here is called a **chart**, which maps neighbourhoods of the manifold to  $\mathbb{R}^k$  (think co-ordinate system or segments of maps). A collection of charts that spans the whole of *M* is called an **atlas**.

**Example** An open subset  $U \subset \mathbb{R}$  is a  $C^{\infty}$  *n*-manifold where the chart is id :  $U \rightarrow U$ .

Notice then in the previous proof,  $\phi$  is a chart, and  $W \cap \{(x,y) \in \mathbb{R}^n : f(x,y) =$ c} is a k-manifold.

For a slightly less trivial example, consider the **unit** n-**sphere**  $S^n = \{x \in \mathbb{R}^{n+1} : \|x\| = 1\}$ . With  $f : \mathbb{R}^{n+1} \to \mathbb{R}$  with  $x \mapsto \|x\|^2$ , we have  $S^n = f^{-1}(\{1\})$ . For every  $x \neq 0$ ,  $Df(x) \neq 0$ , so by implicit function theorem with respect to some co-ordinate system, there exists charts (it turns out an atlas for  $S^n$  requires strictly more than 1 chart). Since f is a polynomial (e.g. standard Cartesian co-ordinates),  $S^n$  is a  $C^\infty$  n-manifold.

Let  $f: U \subset \mathbb{R}^n \to \mathbb{R}^k$  be a  $C^1$  function. A point  $x \in U$  is called a **critical point** if  $\operatorname{rank}(Df(x)) < k$ , i.e. the columns of the derivative matrix do not span  $\mathbb{R}^k$ , and f(x) is called a **critical value**. Otherwise x is called a **regular point**.

**Example** • For  $f : \mathbb{R}^n \to \mathbb{R}$  with f(x) = ||x||, clearly x = 0 is the only critical point, and o is the associated critical value.

- For  $f: \mathbb{R}^n \to \mathbb{R}^k$ , if k > n then there are no regular points in  $\mathbb{R}^n$  by definition.
- For  $f: \mathbb{R}^3 \to \mathbb{R}^2$ , is we have  $f(x, y, z) = (e^z x, (y-1) \sin z)$ , then

$$Df(x,y,z) = \begin{pmatrix} e^z & 0 & e^z x \\ 0 & \sin z & (y-1)\cos z \end{pmatrix}.$$

If  $\sin z \neq 0$  then all points are regular since  $\mathrm{e}^z \neq 0$ . If  $\sin z = 0$ , then  $\cos z = \pm 1$ , and points with  $y \neq 1$  are regular points. Otherwise, the critical points are  $(x,1,n\pi)$  with  $n \in \mathbb{Z}$ , and the critical values are  $f(x,1,n\pi) = (x\mathrm{e}^{n\pi},0)$  (or just the whole y=0 line in  $\mathbb{R}^2$ ).

**Theorem 2.2.6** Let  $f: U \subset \mathbb{R}^n \to \mathbb{R}^k$  be a  $C^{\ell}$ -map with  $\ell \geq 1$ , and U is open. If  $y \in \mathbb{R}^k$  is a regular value, then  $f^{-1}(\{y\})$  is a  $C^{\ell}(n-k)$ -manifold.

**Proof** Let  $x \in f^{-1}(\{y\})$ . Since x is not a critical point, Df(x) has rank k. After rearranging co-ordinates, we can assume that  $(\partial f_i/\partial x_j)(x)$  is invertible, with  $i=1,\ldots k$  and j=n-k+1. The existence of the chart follows from the implicit function theorem, and so  $f^{-1}(\{y\})$  is a  $C^{\ell}(n-k)$ -manifold by definition.

Note that if  $y \notin f(U)$  then  $\phi = f^{-1}(\{y\})$  is still a (n - k)-manifold.

**Example** • For  $f : \mathbb{R}^n \to \mathbb{R}$  with  $x \mapsto ||x||$ , we have  $S^{n-1} = f^{-1}(\{1\})$  following from previous example.

- For  $f(x,y,z)=(\mathrm{e}^z x,(y-1)\sin z)$ , the inverse of the regular values  $f^{-1}(\{(a,b):b\neq 0\})$  is a 1-manifold.
- For  $f: \mathbb{R}^n \times \mathbb{R}^m \to \mathbb{R}^2$  with  $(x,y) \mapsto (\|x\|, \|y\|)$ , we have  $f^{-1}(\{1,1\}) = S^{n-1} \times S^{m-1}$ .

Note that 2-sphere would be the standard sphere, which is two-dimensional with zero volume.

Let  $M \subset \mathbb{R}^n$  be a  $C^\ell$  k-manifold with  $\ell \geq 1$ . The **tangent vector** v at  $p \in M$  is an element in  $\mathbb{R}^n$  of the form  $v = \gamma'(0)$  with  $\gamma : (-\epsilon, \epsilon) \to M$  being a  $C^1$  curve, and that  $\gamma(0) = p$ .

The set of all tangent vectors at point  $p \in M$  is the **tangent space**  $T_p(M)$  at p.

**Proposition 3.0.1** Let  $M \subset \mathbb{R}^n$  be a  $C^\ell$  k-manifold, and  $p \in M$ . Then  $T_p(M)$  is a k-vector space of  $\mathbb{R}^n$ . In fact, if  $h: U \to U' \subset \mathbb{R}^k \times \mathbb{R}^{n-k}$  is a chart with h(p) = 0, then  $T_p(M) \subseteq (Dh^{-1}(0))(\mathbb{R}^k \times \{0\})$ .

**Proof** Let h be a chart,  $\gamma:(-\epsilon,\epsilon)\to M$  with  $\gamma(0)=p\in U$ . We can assume  $\gamma:(-\epsilon,\epsilon)\to U\cap M$ , so

$$h \circ \mathbf{g} : (-\epsilon, \epsilon) \to \mathbb{R}^k \times \{0\},$$

which implies that

$$\gamma = h^{-1} \circ h \circ \gamma,$$

so that

$$v = \gamma'(0) = (Dh^{-1}(h \circ \gamma(0)))(h \circ \gamma)'(0) = Dh^{-1}(0) \cdot w,$$

and thus  $T_p(M) \subseteq Dh^{-1}(0)(\mathbb{R}^k \times \{0\})$ . On the other hand, let  $\delta(t) = tw$ ,  $w \in \mathbb{R}^k$ , and we get a curve in M via  $h^{-1} \circ \delta$ . By the chain rule,

$$(h^{-1} \circ \delta)'(0) = Dh^{-1}(0) \cdot w,$$

which implies that  $Dh^{-1}(0)(\mathbb{R}^k \times \{0\}) \subseteq T_p(M)$ , and so  $Dh^{-1}(0)(\mathbb{R}^k \times \{0\}) = T_p(M)$ .

The chart h is a diffeomorphism so h is injective, which means  $\dim(T_p(M)) = \dim(\mathbb{R}^k) = k$ , as required.

**Theorem 3.0.2** Let  $g: U \subset \mathbb{R}^n \to \mathbb{R}^{n-k}$  be a  $C^{\ell}$ -function, U is open, and  $c \in \mathbb{R}^{n-k}$  is a regular value. Then  $M = g^{-1}(\{c\})$  is a k-manifold and  $T_p(M) = ker\{Dg(p): p \in M\}$ .

Here the kernel is the one induced by the matrix representing the linear map.

**Example** Let  $M = \{(x,y,z) : x^3 + y^3 + z^3 = 1\}$ , and  $g(x,y,z) = x^3 + y^3 + z^3$ . If p = (1,-1,1), then  $T_p(M) = \ker(3(1)^2, 3(-1)^2, 3(1)^2) = \ker(3,3,3) = \{(x,y,z) : x+y+z=0\}$ . On the other hand, for q = (1,0,0), we have  $T_q(M) = \ker(3,0,0) = \{(x,y,z) : x=0\}$ .

Let  $M \subset \mathbb{R}^n$  be a  $C^1$  manifold,  $u \subset \mathbb{R}^n$  open, and  $M \subset U$  with  $f: U \to \mathbb{R}$  a  $C^1$ -function. The point  $p \in M$  is a **critical point** of  $f|_M$  if for every  $C^1$  curve  $\gamma: (-\epsilon, \epsilon) \to M$  with  $\gamma(0) = p$ ,  $(f \circ \gamma)'(0) = 0$ , i.e., the tangent vector is zero at the critical point p.

If  $f|_M$  has a local extreme at  $p \in M$  then p is a critical point. By the chain rule,  $f|_M$  has a critical point exactly when  $Df(p)|_{T_n(M)} = 0$ .

## Method of Lagrange multipliers

**Proposition 3.1.1** Let  $U \subset \mathbb{R}^{n+m}$  be open,  $g: U \to \mathbb{R}^n$  be a  $C^\ell$ -function with  $\ell \geq 1$ , and  $0 \in \mathbb{R}^n$  be a regular value of g. For  $f: U \to \mathbb{R}$  a  $C^1$ -function,  $p \in M = g^{-1}(\{0\})$  is a critical point iff there exists some Lagrange multipliers  $\lambda_1, \ldots \lambda_n \in \mathbb{R}$  with  $D(f + \lambda_i g_i)(p) = 0$ .

**Proof** Assume there exists the relevant Lagrange multipliers, then

$$0 = D(f + \lambda_i g_i)(p) \Leftrightarrow Df(p) = -\lambda_i Dg_i(p).$$

Hence Df(p) is a linear combination of row vectors of  $Dg_i(p)$ . Note that  $Dg_i(p)|_{T_n(M)} = 0$  by the previous theorem, so p is a critical point.

On the other hand, note that  $\operatorname{rank}(Dg(p) = n)$  if p is regular, so  $Dg_i(p)$  are linear independent row vectors. Note also that Df(p) is a linear map from  $\mathbb{R}^{n+m}$  to  $\mathbb{R}$ , vanishing on  $T_p(M)$  which is m-dimensional and sits in the n-dimensional subvector space of the dual space  $(R^{n+m})^*$  housing all of the  $Dg_i(p)$ . Since  $Dg_i(p)$  form a basis for this subspace, we must have constants where  $Df(p) = -\lambda_i Dg_i(p)$ .

The method of Lagrange multipliers gives a method of finding critical points and extrema. Let  $F: U \times \mathbb{R}^n \to \mathbb{R}$  with  $(x, \lambda_1, \dots \lambda_n) \mapsto f(x) + \lambda_i g_i(x)$ , the previous identity gives

$$\frac{\partial F}{\partial x_i} = 0, \qquad \frac{\partial F}{\partial \lambda_j} = 0, \qquad i = 1, \dots n + m, \quad j = 1, \dots n.$$
 (3.1)

Solving the system gives finitely many critical points. Furthermore, if M is compact, then we can find extrema of f via this method.

**Example** Find the maximum value of f(x,y) = x + y on  $M = \{(x,y) : x^4 + y^4 = 1\}.$ 

Defining  $g(x,y) = x^4 + y^4 - 1$ , we have  $g^{-1}(\{0\}) = M$  and is a manifold. We define

$$F(x, y) = f + \lambda_i g_i = x + y + \lambda (x^4 + y^4 - 1),$$

Einstein summation convention implied.

So it is used a lot in optimisation procedures.

which results in

$$0 = \frac{\partial F}{\partial x} = 1 + 4\lambda x^{3},$$
  

$$0 = \frac{\partial F}{\partial y} = 1 + 4\lambda y^{3},$$
  

$$0 = \frac{\partial F}{\partial \lambda} = x^{4} + y^{4} - 1.$$

Since  $(0,0) \notin M$ , the first two equations give

$$x = y = \left(-\frac{1}{4\lambda}\right)^{1/3},$$

so the constraint results in  $\lambda=\pm 8^{1/4}/4$ , and the critical points are  $\pm (2^{-1/4},2^{-1/4})$ . The maximum is thus

$$f(2^{-1/4}, 2^{-1/4}) = \frac{2}{\sqrt{2}}.$$

**Example** Find the extrema of f(x,y,z) = 5x + y - 3z on the intersection of x + y + z = 0 with  $S^2 = \{(x,y,z) : x^2 + y^2 + z^2 = 1\}$ . Consider

$$F(x, y, z, \lambda, \mu) = 5x + y - 3z + \lambda(x + y + z) + \mu(x^2 + y^2 + z^2 - 1).$$

It can be shown that  $\lambda = -1$  from the first three equations. That results in  $y\mu = 0$  in the second equation, and for a non-trivial constraint, we thus have y = 0. This leads then in  $x = -2\mu$ ,  $z = 2/\mu$ , resulting in  $2x^2 = 1$ , and thus the critical points are

$$a = \left(\frac{1}{\sqrt{2}}, 0, -\frac{1}{\sqrt{2}}\right), \quad b = \left(-\frac{1}{\sqrt{2}}, 0, \frac{1}{\sqrt{2}}\right).$$

The extrema are then  $f(a) = 8/\sqrt{2}$  and  $f(b) = -8/\sqrt{2}$ .

**Proposition 3.1.2** Let  $A \subset \mathbb{R}^n$  be compact,  $B \subset \mathbb{R}^n$  be closed, and both non-empty. Then there exists  $a \in A$  and  $b \in B$  where

$$||a-b|| \le ||x-y||$$

for all  $x \in A$  and  $y \in B$ , and this can be any norm.

**Proof** Let  $d = \inf\{||x - y|| : x \in A, y \in B\}$ . For all  $n \in \mathbb{N}$ , there exists some  $a_n \in A$  and  $b_n \in B$  such that

$$||a_n-b_n|| < d+\frac{1}{n}.$$

By passing to a sub-sequence, we can assume  $a_n \to a$  since A is compact. Then we see that

$$||b_n|| \le ||b_n - a_n|| + ||a_n - a|| + ||a|| \le d + 1 + ||a||$$

for  $n \gg 1$ . This implies that  $B \cap D(0; d+1+\|a\|)$  is compact, so  $b_n \to b$  as  $n \to \infty$ . Since  $b \in B$ , we have

$$||a-b|| \le ||a-a_n|| + ||a_n-b_n|| + ||b_n-b|| < d+\epsilon$$

for some  $\epsilon$ . Since d is the infimum, we must have  $||a-b|| \le ||x+y||$  for all  $x \in A$  and  $y \in B$ .

**Example** Find  $q \in M = \{x \in \mathbb{R}^3 : 2x^+y^2 + z = 1\}$  which has minimum distance to p = (0, 0, -5).

Now,  $M = g^{-1}(\{0\})$  where  $g = 2x^2 + y^2 + z - 1$ , and since 0 is a regular value, M is closed (but not bounded). Let  $f(x,y,z) = x^2 + y^2 + (z+5)^2 = ||x-p||^2$  be the norm of choice, and minimising the norm gives us the desired solution. Consider

$$F(x, y, z, \lambda) = x^2 + y^2 + (z+5)^2 + \lambda(2x^2 + y^2 + z - 1).$$

The usual manoeuver gives x = 0 or  $\lambda = -1/2$ , which we consider separately.

- For  $\lambda = -1/2$ , we have y = 0, z = -19/4,  $x = \pm \sqrt{23/8}$ , so  $f(\pm \sqrt{23/8}, 0, -19/4) = 47/16 < 3$ .
- For x = 0, we have y = 0 or  $\lambda = -1$ . The former case gives z = 1 and thus f(0,0,1) = 36 > 3. For  $\lambda = -1$ , we have z = -9/2 and thus  $y = \pm \sqrt{11/2}$ , which gives  $f(0, \pm \sqrt{11/2}, -9/2) = 23/4 > 3$ .

So  $q = (\pm \sqrt{23/8}, 0, -19/4)$ .

# 3.2 Tangent spaces

Let  $M \subset \mathbb{R}^m$  and  $N \subset \mathbb{R}^n$  be two  $C^\ell$  manifolds,  $\ell \geq 1$ . Assume we have a continuous map  $f: M \to N$  which extends to a  $C^1$  map  $\overline{f}: U \to \mathbb{R}^n$ , where  $U \supset M$  is open. We define, for  $p \in M$ ,

$$T_p(\overline{f}): T_p(M) \to T_{\overline{f}(p)}(N),$$

where for  $\gamma:(-\epsilon,\epsilon)\to M$  a  $C^1$  curve with  $\gamma(0)=p$  and  $(\overline{f}\circ\gamma)(0)=\overline{f}(p)$ , we have

$$T_p(\overline{f}) = (f \circ \gamma)'(0) \in T_{f(p)}(N).$$

By chain rule,

$$(\overline{f} \circ \gamma)' = D\overline{f}(\gamma(0)) \cdot \gamma'(0) = D\overline{f}(p) \cdot \gamma'(0),$$

which implies that for  $T_p(\overline{f}): T_p(M) \to T_{\overline{f}(p)}(N)$ , we have

$$T_p(f(v)) = Df(p) \cdot v.$$

We observe that  $v \in T_p(M)$  implies that  $D\overline{f}(p) \cdot v \in T_{\overline{f}(p)}(N)$ , and  $T_p(\overline{f})$  is a linear map between the two tangent spaces.

A map  $f: M \to N$  is called a  $C^{\ell}$  map if f extends to a  $C^{\ell}$  map  $\overline{f}: U \to \mathbb{R}^n$  as before. Here, if  $T_p(f)$  is not surjective, then  $p \in M$  is a **critical point**, and f(p) its **critical value**.

**Theorem 3.2.1** Let  $M \subset \mathbb{R}^m$  and  $N \subset \mathbb{R}^n$  be two  $C^\ell$  manifolds,  $\ell \geq 1$ , and  $f: M \to N$  which extends to a  $C^\ell$  map. If  $x \in N$  is a regular value, then  $f^{-1}(\{x\})$  is a  $C^\ell$  manifold of dimension  $\dim(M) - \dim(N)$ .

**Proof** Let  $y \in f^{-1}(\{x\}) \subseteq M$ , and we seek a chart around y. Let  $s = \dim(M)$  and  $r = \dim(N)$ . We need

$$\psi: U^r \subset \mathbb{R}^{s-r} \times \mathbb{R}^{m+r-s}$$

with

$$\psi\left(f^{-1}(\{x\}\cap U)\right) = U'\cap \left(\mathbb{R}^{s-r}\times\{0\}\right).$$

Let  $g: V \to \mathbb{R}^N \times \mathbb{R}^{n-r}$  be a chart around  $x \in N$ . Choose  $U_y$  such that  $y \in U_y$ ,  $f(U_y) \subset V$ , and a chart

$$h: U_y \to U_y' \subset \mathbb{R}^s \times \mathbb{R}^{m-s}$$
.

We have the following:

$$\mathbb{R}^{s} \times \{0\} \xrightarrow{h} M \xrightarrow{f} N \xrightarrow{\hat{p}} \mathbb{R}^{r} \times \{0\}$$

$$\mathbb{R}^{s} \times \mathbb{R}^{m-s} \qquad \mathbb{R}^{r} \times \mathbb{R}^{n-r}$$

Let  $\phi : \mathbb{R}^s \to \mathbb{R}^r$ , which has full rank since

- a chart maps tangent plane to tangent plane
- *f* is surjective by assumption since *x* is a regular point
- a chart maps to tangent plane.

For  $0 \in \mathbb{R}^s$  corresponding to  $g \in M$  via h,  $D\phi(0)$  has full rank. The same conclusion follows with  $0 \in \mathbb{R}^r$  corresponding to  $x \in N$  via g. Thus  $\phi^{-1}(\{0\})$  is a manifold and corresponds to  $f^{-1}(\{x\}) \cap U_y$  by

$$h\left(\phi^{-1}(\{0\}\times\{0\})\right) = f^{-1}(\{x\})\cap U_y,$$

so  $\phi$  is a chart and  $f^{-1}(\{x\})$  is a manifold.

# 3.3 Vector fields

Let  $M \subset \mathbb{R}^n$  be a  $C^\ell$  manifold with  $\ell \geq 1$ . A continuous function  $v: M \to \mathbb{R}^n$  is called a **vector field** if  $v(x) \in T_x(M)$  for all  $x \in M$ . It is called a  $C^\ell$ -vector field if there is an open set  $U \subset \mathbb{R}^n$  containing M such that v extends to a  $C^\ell$  function  $\overline{v}: U \to \mathbb{R}^n$ .

**Example** For  $S^{n-1} = \{x \in \mathbb{R}^n : ||x||^2 = 1\}$ . Let  $g(x) = ||x||^2$ , then  $S^{n-1} = g^{-1}(\{1\})$ . By theorem,

$$T_p(S^{n-1}) = \text{ker}Dg(p) = \text{ker}(2p) = \{x \in \mathbb{R}^n : 2x_i p_i = 0\}$$
  
=  $\{x \in \mathbb{R}^n : (x, p) = 0\}.$ 

So for a vector field v(x) on  $S^{n-1}$ , we need (x, v(x)) = 0 for all  $x \in S^{n-1}$ . For n = 2m, let  $v : \mathbb{R}^{2m} \to \mathbb{R}^{2m}$ , with

$$(x_1,\ldots,x_{2m})\mapsto (-x_2x_1,-x_3x_2,\ldots,-x_{2m}x_{2m-1}).$$

Here we have (x, v(x)) = 0 for all  $x \in \mathbb{R}^{2m}$ , so v restricts to a vector field on  $S^{2m-1}$ , is  $C^{\infty}$ , and  $v(x) \neq 0$  for all  $x \in S^{2m-1}$ .

v is called a **non-vanishing vector field** in this case. Note that there are no non-vanishing vector fields on  $S^{2m}$ .

**Example** For  $0 < \epsilon < 1$ , define

$$\phi: (1-\epsilon, 1+\epsilon) \times \mathbb{R}^2 \to \mathbb{R}^3; \quad \Phi(r, \phi, \theta) = \begin{pmatrix} (2+r\cos\phi)\cos\theta \\ (2+r\cos\phi)\sin\theta \\ r\sin\phi \end{pmatrix},$$

where  $\phi$  and  $\theta$  are both full angles from 0 to  $2\pi$ . The 2-torus is then

$$T^2 = \{ x \in \mathbb{R}^3 : (x, y, z) = \Phi(1, \phi, \theta) \}.$$

If we restrict  $\Phi$  to small angles we get charts, and so we get tangent plans and vector fields. Note that the 2-torus does have non-vanishing vector fields, compared to the 2-sphere.

Let  $U, V \subset \mathbb{R}^n$  be open sets,  $h: U \times V$  be a  $C^{\infty}$  diffeomorphism, and  $v: U \to \mathbb{R}^n$  be a vector field. We define the vector field on v by

$$h * v : V \to \mathbb{R}^n$$
,  $h * v(x) = Dh\left(h^{-1}(x)\right) \cdot v(h^{-1}(x))$ .

**Lemma 3.3.1** If M is a  $C^{\infty}$  manifold and v a  $C^{\ell}$  vector field, with  $\ell \geq 1$ . For all  $p \in M$ , there exists open  $I \subset \mathbb{R}$  with  $0 \in I$ , and an integral curve  $\gamma: I \to M$  such that  $\gamma(0) = p$ ,  $\gamma'(t) = v(\gamma(t))$  for all  $t \in I$ .

**Lemma 3.3.2** As above, for i = 1, 2, let  $\gamma_i : I_i \to M$  be integral curves of v with  $\gamma_1(0) = p = \gamma_2(0)$ ,  $I_i$  open, and  $0 \in I_i$ . Then  $\gamma_1(t) = \gamma_2(t)$  for all  $t \in I_1 \cap I_2$ .

**Proof** Uniqueness follows from the Picard–Lindelöf theorem.

Note that the integral curve can now be extended to an integral curve of  $I_1 \cup I_2$ , and we get a maximal curve through a point p this way.

This is related to the **hairy ball theorem**.

**Proposition 3.3.3** Let M be a compact  $C^{\infty}$  manifold and v and  $C^{\ell}$  vector *field with*  $\ell \geq 1$ *. For all*  $p \in M$ *, there exists*  $\gamma : \mathbb{R} \to M$  *with*  $\gamma(0) = p$ *.* 

**Proof** Let  $\gamma: I \to M$  be the maximal integral curve, and assume  $I \cap [0, \infty)$  is bounded. Then these exists  $T = \sup\{I \cap [0, \infty)\}$ . Choose a sequence  $(t_n) \in I$  with  $t_n \to T$ , then  $\gamma(t_n)$  is a sequence in M. Since M is compact, we can assume  $\gamma(t_n) \to x \in M$ .

Let  $\beta$  :  $(T - \epsilon, T + \epsilon) \rightarrow M$  be an integral curve with  $\beta(T) = x$ . Since  $t_n \to T$  for large n, and  $t_n \in (T - \epsilon, T + \epsilon)$ , we should have  $\gamma(t_n) = \beta(t_n)$  by uniqueness, and so  $\gamma$  can be extended beyond T. However, this is a contradiction since  $\gamma$  was assumed to be maximal, so *I* is not bounded, and thus  $\gamma$  can be extended to  $\mathbb{R}$ .

For v a  $C^{\ell}$  vector field ( $\ell \geq 1$ ) on a compact manifold M, the **flow**  $\Phi$  is defined as

$$\Phi: M \times \mathbb{R} \to M, \qquad (x,t) \mapsto \gamma_x(t)$$
 (3.2)

where  $\gamma_x$  is the integral curve with  $\gamma_x(0) = x$ .

**Theorem 3.3.4** Let M be a compact  $C^{\infty}$  manifold and v a  $C^{\ell}$  vector field,  $\ell \geq 1$ . Then the flow  $\Phi$  is continuous and

1.  $\Phi(x,0) = x$  for all  $x \in M$ ,

2. 
$$\Phi(\Phi(x,t),x) = \Phi(x,t+s)$$
 for all  $x \in M$ ,  $t,s \in \mathbb{R}$ .

Proof Continuity holds and follows from Picard-Lindelöf, and  $\Phi(x,0) = x$  follows from definition of the flow map. Let  $y = \Phi(x,t)$ , so  $\gamma_x(t) = y$ . Define  $\gamma(u) = \gamma_x(u+t)$ , which is an integral curve with  $\gamma(0) = y$ . By uniqueness,  $\gamma = \gamma_y$ , and so

$$\Phi(\Phi(x,t),s) = \gamma_y(s) = \gamma(s) = \gamma_x(s+t) = \Phi(x,t+s).$$

Note that if we write  $x \cdot t = \Phi(x, t)$ , then  $x \cdot 0 = x$ , and  $(x \cdot t) \cdot t$  $s = x \cdot (t + s)$ , so the abelian group  $\mathbb{R}$  acts on the set M. Since  $\Phi$  is continuous we have a topological action. Every  $C^1$  vector field v on a compact manifold M gives rise to an  $\mathbb{R}$ -action on M.

Note also that v is of  $C^{\ell}$  class implies that  $\Phi$  is of  $C^{\ell}$  class.

In lower dimensions, from the standard fundamental theorem of calculus, Stokes' theorem and divergence theorem, we see we have identities of the form

$$\int_{M} d\omega = \int_{\partial M} \omega, \tag{4.1}$$

where M is some (oriented) manifold, and  $\omega$  is some function / vector field. This is in fact true in higher dimensions, and the result is the **generalised Stokes' theorem**. It will be seen  $\omega$  is a **differential** k-**form**, and M are the **oriented**  $\ell$ -**manifolds** in  $\mathbb{R}^n$  with boundary  $\partial M$ . To get to the general result, we go through some machinery first in  $\mathbb{R}^n$ , before proceeding to general (oriented) manifolds.

# 4.1 Riemann integrals

For  $f : [a, b] \to \mathbb{R}$ , recall that for a partition  $Z = \{t_0, t_1, \dots, t_n\}$ , the **upper/lower Riemann sums** are defined as

$$\mathcal{U}(f,Z) = \sum_{i=0}^{n-1} M_i(f)(t_{i+1} - t_i), \qquad \mathcal{L}(f,Z) = \sum_{i=0}^{n-1} m_i(f)(t_{i+1} - t_i),$$
(4.2)

where for  $x \in [t_{i-1}, t_i]$ ,

$$M_i(f) = \sup f(x), \quad m_i(f) = \inf f(x).$$

If Z' is a **refinement** of Z (i.e.  $Z' \supset Z$ , where Z' is a partition), then

$$\mathcal{L}(f,Z) \leq \mathcal{L}(f,Z') \leq \mathcal{U}(f,Z') \leq \mathcal{U}(f,Z).$$

For two partitions, the **common refinement** is  $Z'' = Z' \cup Z$ , which implies that

$$\mathcal{L}(f, Z) \le \mathcal{L}(f, Z'') \le \mathcal{U}(f, Z'') \le \mathcal{U}(f, Z').$$

The upper Riemann integral is then defined as

$$\int_{[a,b]}^{u} f \, \mathrm{d}x = \inf \{ \mathcal{U}(f,Z) : Z \text{ a partition of } [a,b] \}, \tag{4.3}$$

while the lower Riemann integral is

$$\int_{[a,b]}^{l} f \, \mathrm{d}x = \inf\{\mathcal{L}(f,Z) : Z \text{ a partition of } [a,b]\}. \tag{4.4}$$

For bounded f, we should have

$$\int_{[a,b]}^{l} f \, \mathrm{d}x \le \int_{[a,b]}^{u} f \, \mathrm{d}x \le \infty.$$

If the two sums coincide as  $|t_{i-1} - t_i| \rightarrow 0$ , then f is **Riemann** integrable.

In  $\mathbb{R}^n$ , to generalise, f defined analogously if each individual component of f is Riemann integrable.

**Lemma 4.1.1** *Let*  $f:[a,b] \to \mathbb{R}^n$  *be integrable. Then* ||f|| *is also integrable and* 

$$\left\| \int_{[a,b]} f \, \mathrm{d}x \right\| \le \int_{[a,b]} \|f\| \, \mathrm{d}x.$$

**Proof** ||f|| is clear integrable. We see that for all  $\epsilon > 0$ , there exists a common partition Z such that

$$\mathcal{U}(f_i, Z) - \mathcal{L}(f_i, Z) \le \epsilon, \quad \mathcal{U}(\|f\|, Z) - \mathcal{L}(\|f\|, Z) \le \epsilon$$

for all components  $f_i$  of f. For any partition  $Z = \{x_0, \dots, x_n\}$  and any choice  $\xi_i$ 

$$a = x_0 \le \xi_0 \le x_1 \le \xi_1 \le \dots \le x_{n-1} \le \xi_{n-1} \le x_n = b$$

we have

$$\left\| \sum_{i=0}^{n-1} f(\xi_i)(x_{i+1} - x_i) \right\| \le \sum_{i=0}^{n-1} \| f(\xi_i) \| (x_{i+1} - x_i)$$

by the triangle inequality. For such a partition Z, we have both

$$\left| \int_{a}^{b} f_{i} \, dx - \sum_{k=0}^{n-1} f_{i}(\xi_{k})(x_{k+1} - x_{k}) \right| \leq \epsilon,$$

$$\left| \int_{a}^{b} \|f\| \, dx - \sum_{k=0}^{n-1} \|f(\xi_{k})\|(x_{k+1} - x_{k}) \right| \leq \epsilon,$$

which implies that, considering each component,

$$\left\| \int_a^b f \, \mathrm{d}x - \sum_{k=0}^{n-1} f(\xi_k) (x_{k+1} - x_k) \right\| \le \sqrt{n\epsilon^2} = \sqrt{n\epsilon}.$$

Then,

$$\left\| \int_{[a,b]} f \, dx \right\| \le \left\| \sum_{k=0}^{n-1} f(\xi_k) (x_{k+1} - x_k) \right\| + \sqrt{n} \epsilon$$

$$\le \sum_{k=0}^{n-1} \| f(\xi_k) (x_{k+1} - x_k) \| + \sqrt{n} \epsilon$$

$$\le \int_a^b \| f \| \, dx + \epsilon + \sqrt{n} \epsilon.$$

Since  $\epsilon$  was arbitrary, we have the result as required.

# 4.2 Differential 1-forms and line integrals

Let *V* be a real vector space with norm  $\|\cdot\|$ , and  $c:[a,b]\to V$  be continuous. The **length** of *c* is defined as

$$L(c) = \sup \left\{ \sum_{i=1}^{n-1} \|c(t_{i+1} - c(t_i))\| : \forall n \in N, t_i \text{ in } Z \right\}.$$

The curve *c* is **rectifiable** if  $L(c) < \infty$ .

Note that the length of a curve is (and should be) independent of its parameterisation.

**Proposition 4.2.1** For  $c:[a,b] \to^n$  of class  $C^1$ , c is rectifiable, and  $L(c) = \int_a^b \|c'(t)\| dt$ .

**Proof** Note that

$$\sum_{i} \|c(t_{i+1} - c(t_{i})\| = \sum_{i} \left\| \int_{[t_{i}, t_{i+1}]} c'(t) \, dt \right\|$$

$$\leq \sum_{i} \int_{[t_{i}, t_{i+1}]} \|c'(t)\| \, dt$$

$$= \int_{a}^{b} \|c'(t)\| \, dt < \infty$$

since  $c \in C^1[a, b]$ , so  $L(c) < \infty$  and c is rectifiable. Let  $f(t) = L\left(c|_{[a,t]}\right)$  for  $a \le t_0 < t \le b$ . Then

$$\left| \frac{c(t) - c(t_0)}{t - t_0} \right| \le \frac{L\left(c|_{[t_0, t]}\right)}{t - t_0} = \frac{f(t) - f(t_0)}{t - t_0}$$

$$\le \frac{1}{t - t_0} \int_{[t_0, t]} \|c'(s)\| \, \mathrm{d}s = \|c'(t_1)\|$$

for  $t_1 \in [t_0, t]$  by the mean value theorem. So as  $t_0, t_1 \to t$ ,

$$||c'(t)|| \le f'(t) \le ||c'(t)||,$$

and so f'(t) exists and f'(t) = ||c'(t)||. Thus

$$L(c) = f(b) - f(a) = \int_a^b f'(t) dt = \int_a^b ||c'(t)|| dt < \infty$$

as required.

**Example** For the helix, we have  $c(t) = (at, r \cos t, r \sin t)$ , so  $||c'(t)|| = \sqrt{r^2 + a^2}$ , and

$$L(c)|_{[0,2\pi]} = \int_0^{2\pi} \sqrt{r^2 + a^2} \, dt = 2\pi \sqrt{r^2 + a^2}.$$

Recall the differential  $Df(x): \mathbb{R}^n \to \mathbb{R}$  is a linear map. Denote this now has  $df_x : \mathbb{R}^n \to \mathbb{R}$ , and the real vector space of all linear maps  $\phi: \mathbb{R}^n \to \mathbb{R}$  to be  $\mathcal{L}(\mathbb{R}^n, \mathbb{R})$ . Since  $\phi \in \mathcal{L}(\mathbb{R}^n, \mathbb{R})$  is linear, we only investigate the action of  $\phi$  on the basis  $\{e_i\}$ . We see that

$$\mathrm{d}f_x(e_i) = \frac{\partial f}{\partial x_i}(x),$$

so

$$\mathrm{d}f_x(v) = \frac{\partial f}{\partial x_i} v_i = \langle \nabla f, v \rangle$$

for  $\langle \cdot, \cdot \rangle$  the inner product on some vector space V housing v. We see  $\partial f/\partial x_i$  are smooth coefficient functions, and that since  $f:U\to\mathbb{R}^n$  is smooth,  $df: U \to \mathcal{L}(\mathbb{R}^n, \mathbb{R})$ .

Let  $U \subset \mathbb{R}^n$  be open and  $\omega : U \to \mathcal{L}(\mathbb{R}^n, \mathbb{R})$ , then there are functions  $f_1, \dots f_n : U \to \mathbb{R}^n$  where

$$\omega_x(v) = f_i(x)v_i \tag{4.5}$$

for all  $x \in U$ ,  $v = v_i e_i$ . The coefficient functions  $f_i$  are calculated via

$$f_i(x) = \omega_x(e_i). \tag{4.6}$$

We call  $\omega$  a **differential 1-form** on U if  $f_i:U\to\mathbb{R}$  are of class  $C^{\infty}$  for all i. The set of all 1-forms is denoted by  $\Omega^{1}(U)$ , which has the structure of a real vector space. One can canonically multiply a 1-form  $\omega \in \Omega^1(U)$  with a smooth function  $f \in C^{\infty}(U)$  by performing

$$(f\omega)_{x}(v) = f(x)\omega_{x}(v).$$

**Lemma 4.2.2** Let  $\omega$  be a 1-form and  $U \subset \mathbb{R}^N$  be open. Then there exists a smooth vector field  $X_{\omega}: U \to \mathbb{R}^n$  where

$$\omega_x(v) = \langle X_{\omega}(x), v \rangle.$$

**Proof** Since  $F_{\omega} = f_i e_i$ ,  $\langle F_{\omega}(x), v \rangle = f_i(x) \langle e_i, v \rangle = f_i(x) v_i = \omega_x(v)$ .

**Lemma 4.2.3** Let  $\omega \in \Omega^1(U)$ ,  $x_i : U \to \mathbb{R}$ ,  $x_i(p_1, \ldots, p_n) = p_i$ . Then  $\omega = f_i \, \mathrm{d} x_i$  and  $f_i(x) = \omega_x(e_i)$ .

**Proof** Since

$$\mathrm{d}x_i(x)(v) = \frac{\partial x_i}{\partial x_j}(x)v_j = v_i,$$

we have

$$\omega_x(v) = f_i(x)v_i = f_i(x) \, \mathrm{d}x_i(x)(v),$$

and since v was arbitrary,  $\omega_x = f_i(x) dx_i$ .

Einstein summation.

**Example** Suppose  $\omega \in \Omega^1(\mathbb{R}^2)$  and  $\omega = 3xy \, dx + y^3 \, dy$ , and we take p = (7,3). Then since

$$\omega_p(e_1) = 3 \cdot p_1 \cdot p_2 \, dx(e_1) + p_2^3 \, dy(e_1) = 3 \cdot 7 \cdot 1 + 27 \cdot 0 = 63,$$

while  $\omega_p(e_2)=27$  by a similar argument. So  $\omega_p((1,-2))=63-2\cdot 27=9$  for example.

A differential 1-form  $\omega$  is **exact** if there exists some  $f \in C^\infty(U)$  where

$$\omega = \mathrm{d}f. \tag{4.7}$$

For  $U \subset \mathbb{R}^n$  open and  $c : [a,b] \to U$  be smooth and  $\omega \in \Omega^1(U)$ . The **line integral** of  $\omega$  along c is

$$\int_{C} \omega = \int_{a}^{b} \omega_{c(t)} \left( c'(t) \right) dt. \tag{4.8}$$

If *c* is piecewise smooth, we can still define the integral.

For  $c:[a,b]\to\mathbb{R}^n$  a smooth curve and  $\phi:[\alpha,\beta]\to[a,b]$  be a smooth bijective map, then  $\tilde{c}=c\circ\phi:[\alpha,\beta]\to\mathbb{R}^n$  is a **orientation preserving reparameterisation** if  $\phi'>0$ , otherwise it is orientation reversing.

**Proposition 4.2.4** For  $\omega \in \Omega^1(U)$ ,  $c : [a,b] \to U$ ,  $\tilde{c} = c \circ \phi : [\alpha,\beta] \to U$  orientation preserving,<sup>1</sup>

$$\int_{\mathcal{E}} \omega = \int_{\tilde{\mathcal{E}}} \omega.$$

**Proof** 

$$\int_{\tilde{c}} \omega = \int_{\alpha}^{\beta} \omega_{\tilde{c}}(\tilde{c}') dt = \int_{\alpha}^{\beta} \omega_{c \circ \phi} ((c \circ \phi)') dt$$
$$= \int_{\alpha}^{\beta} \omega_{c \circ \phi} ((c' \circ \phi)) \phi' dt$$
$$= \int_{a}^{b} \omega_{c} (c') dt = \int_{c} \omega.$$

**Lemma 4.2.5** For  $f \in C^{\infty}(U)$ ,  $c : [a, b] \rightarrow U$ , then

$$\int_{c} \mathrm{d}f = f(c(b)) - f(c(a)).$$

**Proof** 

$$\int_{c} df = \int_{a}^{b} [Df(c(t))](c'(t)) dt = \int_{a}^{b} (f \circ c)'(t) dt = f(c(b)) - f(c(a)).$$

**Proposition 4.2.6** *Let*  $U \subset \mathbb{R}^n$  *be open and path connected, then the following are equilvalent:* 

A 1-form eats a vector and spits out a number. Sometimes it can be regarded as a functional (eats a function and spits out a number).

<sup>1</sup> If orientation reversing, then there would be an extra minus sign.

- 1.  $\omega \in \Omega^1(U)$  is exact,
- 2.  $\int_{C} \omega$  depends only on the end points (i.e. path independence),
- 3.  $\oint_{\mathcal{C}} \omega = 0$ .

**Proof** • 1) implies 2) by previous lemma.

- 2) implies 3). Consider c with  $c(a)=c(b)=p\in U$ , and  $c_p(t)=p$  for all t. Then since  $c_p'(t)=0$ ,  $\int_{c_p}\omega=0$  and so  $\oint_c\omega=0$  since we assumed path independence.
- 3) implies 2). For  $c = c_1 \cup (-c_2)$ ,

$$0 = \oint_c \omega = \int_{c_1} \omega + \int_{-c_2} \omega = \int_{c_1} \omega - \int_{c_2} \omega,$$

which implies path independence.

• 2) implies 1). Choosing  $p \in U$ ,  $f : U \to \mathbb{R}$  via  $f(q) = \int_c \omega$ , c(0) = p, c(1) = q, then f is well-defined by 2). Choose  $h \in \mathbb{R}^n$  where  $h + q \in U$ , let

$$\gamma:[0,1]\to I,\quad \gamma(t)=q+th.$$

Then

$$f(q+th)-f(q)=\int_{\gamma}\omega=\int_{0}^{1}\omega_{\gamma}(\gamma')\;\mathrm{d}t\int_{0}^{1}\omega_{q+th}(h)\;\mathrm{d}t.$$

Now introduce vector field  $X_{\omega}$ , an by lemma,

$$f(q+th) - f(q) - \omega_q(h) = \int_0^1 \left( \omega_{q+th}(h) - \omega_q(h) \right) dt$$
$$= \int_0^1 \langle X_{\omega}(q+th) - X_{\omega}(q), h \rangle dt$$
$$\leq \int_0^1 \|X_{\omega}(q+th) - X_{\omega}(q)\| dt \cdot \|h\|.$$

The right hand side has  $R(h)/\|h\| \to 0$  as  $h \to 0$  since  $X_{\omega}$  is continuous, so f is differentiable, and

$$\omega_q(h) = Df(q)(h) = \mathrm{d}f_q(h),$$

and since q and h are arbitrary,  $\omega = df$ .

**Proposition 4.2.7** *If*  $\omega \in \Omega^1(U)$  *is exact and*  $\omega = f_i dx_i$ , then

$$\frac{\partial f_i}{\partial x_i} = \frac{\partial f_j}{\partial x_i}$$

for all i and j.

**Proof** Since

$$\omega = \mathrm{d}f = \frac{\partial f}{\partial x_i} \, \mathrm{d}x_i,$$

we have  $f_i = \partial f / \partial x_i$ , so then since f is differentiable,

$$\frac{\partial f_i}{\partial x_j} = \frac{\partial^2 f}{\partial x_j \partial x_i} = \frac{\partial^2 f}{\partial x_i \partial x_j} = \frac{\partial^2 f_j}{\partial x_i}.$$

 $\omega \in \Omega^1(U)$  is **closed** if for  $\omega = f_i \, \mathrm{d} x_i, \, \partial f_i/\partial x_j = \partial f_j/\partial x_i.$ 

- 4.3 *Differential k-forms*
- 4.4 Integration in  $\mathbb{R}^n$

# Differential forms on oriented manifolds

Stokes' theorem