

List of Courses

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**1/I/1F Analysis II**

Let  $E$  be a subset of  $\mathbb{R}^n$ . Prove that the following conditions on  $E$  are equivalent:

(i)  $E$  is closed and bounded.

(ii)  $E$  has the Bolzano–Weierstrass property (i.e., every sequence in  $E$  has a subsequence convergent to a point of  $E$ ).

(iii) Every continuous real-valued function on  $E$  is bounded.

[The Bolzano–Weierstrass property for bounded closed intervals in  $\mathbb{R}^1$  may be assumed.]

**1/II/10F Analysis II**

Explain briefly what is meant by a *metric space*, and by a *Cauchy sequence* in a metric space.

A function  $d : X \times X \rightarrow \mathbb{R}$  is called a pseudometric on  $X$  if it satisfies all the conditions for a metric except the requirement that  $d(x, y) = 0$  implies  $x = y$ . If  $d$  is a pseudometric on  $X$ , show that the binary relation  $R$  on  $X$  defined by  $x R y \Leftrightarrow d(x, y) = 0$  is an equivalence relation, and that the function  $d$  induces a metric on the set  $X/R$  of equivalence classes.

Now let  $(X, d)$  be a metric space. If  $(x_n)$  and  $(y_n)$  are Cauchy sequences in  $X$ , show that the sequence whose  $n$ th term is  $d(x_n, y_n)$  is a Cauchy sequence of real numbers. Deduce that the function  $\bar{d}$  defined by

$$\bar{d}((x_n), (y_n)) = \lim_{n \rightarrow \infty} d(x_n, y_n)$$

is a pseudometric on the set  $C$  of all Cauchy sequences in  $X$ . Show also that there is an isometric embedding (that is, a distance-preserving mapping)  $X \rightarrow C/R$ , where  $R$  is the equivalence relation on  $C$  induced by the pseudometric  $\bar{d}$  as in the previous paragraph. Under what conditions on  $X$  is  $X \rightarrow C/R$  bijective? Justify your answer.

**2/I/1F Analysis II**

Explain what it means for a function  $f : \mathbb{R}^2 \rightarrow \mathbb{R}^1$  to be *differentiable* at a point  $(a, b)$ . Show that if the partial derivatives  $\partial f / \partial x$  and  $\partial f / \partial y$  exist in a neighbourhood of  $(a, b)$  and are continuous at  $(a, b)$  then  $f$  is differentiable at  $(a, b)$ .

Let

$$f(x, y) = \frac{xy}{x^2 + y^2} \quad ((x, y) \neq (0, 0))$$

and  $f(0, 0) = 0$ . Do the partial derivatives of  $f$  exist at  $(0, 0)$ ? Is  $f$  differentiable at  $(0, 0)$ ? Justify your answers.

**2/II/10F Analysis II**

Let  $V$  be the space of  $n \times n$  real matrices. Show that the function

$$N(A) = \sup \{ \|A\mathbf{x}\| : \mathbf{x} \in \mathbb{R}^n, \|\mathbf{x}\| = 1 \}$$

(where  $\| - \|$  denotes the usual Euclidean norm on  $\mathbb{R}^n$ ) defines a norm on  $V$ . Show also that this norm satisfies  $N(AB) \leq N(A)N(B)$  for all  $A$  and  $B$ , and that if  $N(A) < \epsilon$  then all entries of  $A$  have absolute value less than  $\epsilon$ . Deduce that any function  $f: V \rightarrow \mathbb{R}$  such that  $f(A)$  is a polynomial in the entries of  $A$  is continuously differentiable.

Now let  $d: V \rightarrow \mathbb{R}$  be the mapping sending a matrix to its determinant. By considering  $d(I + H)$  as a polynomial in the entries of  $H$ , show that the derivative  $d'(I)$  is the function  $H \mapsto \text{tr } H$ . Deduce that, for any  $A$ ,  $d'(A)$  is the mapping  $H \mapsto \text{tr}((\text{adj } A)H)$ , where  $\text{adj } A$  is the adjugate of  $A$ , i.e. the matrix of its cofactors.

[*Hint: consider first the case when  $A$  is invertible. You may assume the results that the set  $U$  of invertible matrices is open in  $V$  and that its closure is the whole of  $V$ , and the identity  $(\text{adj } A)A = \det A \cdot I$ .]*

**3/I/1F Analysis II**

Let  $V$  be the vector space of continuous real-valued functions on  $[-1, 1]$ . Show that the function

$$\|f\| = \int_{-1}^1 |f(x)| dx$$

defines a norm on  $V$ .

Let  $f_n(x) = x^n$ . Show that  $(f_n)$  is a Cauchy sequence in  $V$ . Is  $(f_n)$  convergent? Justify your answer.

**3/II/11F Analysis II**

State and prove the Contraction Mapping Theorem.

Let  $(X, d)$  be a bounded metric space, and let  $F$  denote the set of all continuous maps  $X \rightarrow X$ . Let  $\rho: F \times F \rightarrow \mathbb{R}$  be the function

$$\rho(f, g) = \sup \{ d(f(x), g(x)) : x \in X \} .$$

Show that  $\rho$  is a metric on  $F$ , and that  $(F, \rho)$  is complete if  $(X, d)$  is complete. [*You may assume that a uniform limit of continuous functions is continuous.*]

Now suppose that  $(X, d)$  is complete. Let  $C \subseteq F$  be the set of contraction mappings, and let  $\theta: C \rightarrow X$  be the function which sends a contraction mapping to its unique fixed point. Show that  $\theta$  is continuous. [*Hint: fix  $f \in C$  and consider  $d(\theta(g), f(\theta(g)))$ , where  $g \in C$  is close to  $f$ .]*

4/I/1F **Analysis II**

Explain what it means for a sequence of functions  $(f_n)$  to converge uniformly to a function  $f$  on an interval. If  $(f_n)$  is a sequence of continuous functions converging uniformly to  $f$  on a finite interval  $[a, b]$ , show that

$$\int_a^b f_n(x) dx \longrightarrow \int_a^b f(x) dx \quad \text{as } n \rightarrow \infty .$$

Let  $f_n(x) = x \exp(-x/n)/n^2$ ,  $x \geq 0$ . Does  $f_n \rightarrow 0$  uniformly on  $[0, \infty)$ ? Does  $\int_0^\infty f_n(x) dx \rightarrow 0$ ? Justify your answers.

4/II/10F **Analysis II**

Let  $(f_n)_{n \geq 1}$  be a sequence of continuous complex-valued functions defined on a set  $E \subseteq \mathbb{C}$ , and converging uniformly on  $E$  to a function  $f$ . Prove that  $f$  is continuous on  $E$ .

State the Weierstrass  $M$ -test for uniform convergence of a series  $\sum_{n=1}^\infty u_n(z)$  of complex-valued functions on a set  $E$ .

Now let  $f(z) = \sum_{n=1}^\infty u_n(z)$ , where

$$u_n(z) = n^{-2} \sec(\pi z/2n) .$$

Prove carefully that  $f$  is continuous on  $\mathbb{C} \setminus \mathbb{Z}$ .

[You may assume the inequality  $|\cos z| \geq |\cos(\operatorname{Re} z)|$ .]

**1/I/7B Complex Methods**

Let  $u(x, y)$  and  $v(x, y)$  be a pair of conjugate harmonic functions in a domain  $D$ . Prove that

$$U(x, y) = e^{-2uv} \cos(u^2 - v^2) \quad \text{and} \quad V(x, y) = e^{-2uv} \sin(u^2 - v^2)$$

also form a pair of conjugate harmonic functions in  $D$ .

**1/II/16B Complex Methods**

Sketch the region  $A$  which is the intersection of the discs

$$D_0 = \{z \in \mathbb{C} : |z| < 1\} \quad \text{and} \quad D_1 = \{z \in \mathbb{C} : |z - (1 + i)| < 1\}.$$

Find a conformal mapping that maps  $A$  onto the right half-plane  $H = \{z \in \mathbb{C} : \operatorname{Re} z > 0\}$ . Also find a conformal mapping that maps  $A$  onto  $D_0$ .

[Hint: You may find it useful to consider maps of the form  $w(z) = \frac{az+b}{cz+d}$ .]

**2/I/7B Complex Methods**

(a) Using the residue theorem, evaluate

$$\int_{|z|=1} \left(z - \frac{1}{z}\right)^{2n} \frac{dz}{z}.$$

(b) Deduce that

$$\int_0^{2\pi} \sin^{2n} t \, dt = \frac{\pi}{2^{2n-1}} \frac{(2n)!}{(n!)^2}.$$

**2/II/16B Complex Methods**

(a) Show that if  $f$  satisfies the equation

$$f''(x) - x^2 f(x) = \mu f(x), \quad x \in \mathbb{R}, \quad (*)$$

where  $\mu$  is a constant, then its Fourier transform  $\widehat{f}$  satisfies the same equation, i.e.

$$\widehat{f}''(\lambda) - \lambda^2 \widehat{f}(\lambda) = \mu \widehat{f}(\lambda).$$

(b) Prove that, for each  $n \geq 0$ , there is a polynomial  $p_n(x)$  of degree  $n$ , unique up to multiplication by a constant, such that

$$f_n(x) = p_n(x)e^{-x^2/2}$$

is a solution of (\*) for some  $\mu = \mu_n$ .

(c) Using the fact that  $g(x) = e^{-x^2/2}$  satisfies  $\widehat{g} = cg$  for some constant  $c$ , show that the Fourier transform of  $f_n$  has the form

$$\widehat{f}_n(\lambda) = q_n(\lambda)e^{-\lambda^2/2}$$

where  $q_n$  is also a polynomial of degree  $n$ .

(d) Deduce that the  $f_n$  are eigenfunctions of the Fourier transform operator, i.e.  $\widehat{f}_n(x) = c_n f_n(x)$  for some constants  $c_n$ .

**4/I/8B Complex Methods**

Find the Laurent series centred on 0 for the function

$$f(z) = \frac{1}{(z-1)(z-2)}$$

in each of the domains

$$(a) \quad |z| < 1, \quad (b) \quad 1 < |z| < 2, \quad (c) \quad |z| > 2.$$

4/II/17B **Complex Methods**

Let

$$f(z) = \frac{z^m}{1+z^n}, \quad n > m+1, \quad m, n \in \mathbb{N},$$

and let  $C_R$  be the boundary of the domain

$$D_R = \{z = re^{i\theta} : 0 < r < R, \quad 0 < \theta < \frac{2\pi}{n}\}, \quad R > 1.$$

(a) Using the residue theorem, determine

$$\int_{C_R} f(z) dz.$$

(b) Show that the integral of  $f(z)$  along the circular part  $\gamma_R$  of  $C_R$  tends to 0 as  $R \rightarrow \infty$ .

(c) Deduce that

$$\int_0^\infty \frac{x^m}{1+x^n} dx = \frac{\pi}{n \sin \frac{\pi(m+1)}{n}}.$$

**1/I/6C Fluid Dynamics**

An unsteady fluid flow has velocity field given in Cartesian coordinates  $(x, y, z)$  by  $\mathbf{u} = (1, xt, 0)$ , where  $t$  denotes time. Dye is released into the fluid from the origin continuously. Find the position at time  $t$  of the dye particle that was released at time  $s$  and hence show that the dye streak lies along the curve

$$y = \frac{1}{2}tx^2 - \frac{1}{6}x^3.$$

**1/II/15C Fluid Dynamics**

Starting from the Euler equations for incompressible, inviscid flow

$$\rho \frac{D\mathbf{u}}{Dt} = -\nabla p, \quad \nabla \cdot \mathbf{u} = 0,$$

derive the vorticity equation governing the evolution of the vorticity  $\boldsymbol{\omega} = \nabla \times \mathbf{u}$ .

Consider the flow

$$\mathbf{u} = \beta(-x, -y, 2z) + \Omega(t)(-y, x, 0),$$

in Cartesian coordinates  $(x, y, z)$ , where  $t$  is time and  $\beta$  is a constant. Compute the vorticity and show that it evolves in time according to

$$\boldsymbol{\omega} = \omega_0 e^{2\beta t} \mathbf{k},$$

where  $\omega_0$  is the initial magnitude of the vorticity and  $\mathbf{k}$  is a unit vector in the  $z$ -direction.

Show that the material curve  $C(t)$  that takes the form

$$x^2 + y^2 = 1 \quad \text{and} \quad z = 1$$

at  $t = 0$  is given later by

$$x^2 + y^2 = a^2(t) \quad \text{and} \quad z = \frac{1}{a^2(t)},$$

where the function  $a(t)$  is to be determined.

Calculate the circulation of  $\mathbf{u}$  around  $C$  and state how this illustrates Kelvin's circulation theorem.



**3/I/8C Fluid Dynamics**

Show that the velocity field

$$\mathbf{u} = \mathbf{U} + \frac{\mathbf{\Gamma} \times \mathbf{r}}{2\pi r^2},$$

where  $\mathbf{U} = (U, 0, 0)$ ,  $\mathbf{\Gamma} = (0, 0, \Gamma)$  and  $\mathbf{r} = (x, y, 0)$  in Cartesian coordinates  $(x, y, z)$ , represents the combination of a uniform flow and the flow due to a line vortex. Define and evaluate the circulation of the vortex.

Show that

$$\oint_{C_R} (\mathbf{u} \cdot \mathbf{n}) \mathbf{u} \, dl \rightarrow \frac{1}{2} \mathbf{\Gamma} \times \mathbf{U} \quad \text{as} \quad R \rightarrow \infty,$$

where  $C_R$  is a circle  $x^2 + y^2 = R^2$ ,  $z = \text{const}$ . Explain how this result is related to the lift force on a two-dimensional aerofoil or other obstacle.

**3/II/18C Fluid Dynamics**

State the form of Bernoulli's theorem appropriate for an unsteady irrotational motion of an inviscid incompressible fluid in the absence of gravity.

Water of density  $\rho$  is driven through a tube of length  $L$  and internal radius  $a$  by the pressure exerted by a spherical, water-filled balloon of radius  $R(t)$  attached to one end of the tube. The balloon maintains the pressure of the water entering the tube at  $2\gamma/R$  in excess of atmospheric pressure, where  $\gamma$  is a constant. It may be assumed that the water exits the tube at atmospheric pressure. Show that

$$R^3 \ddot{R} + 2R^2 \dot{R}^2 = -\frac{\gamma a^2}{2\rho L}. \quad (\dagger)$$

Solve equation  $(\dagger)$ , by multiplying through by  $2R\dot{R}$  or otherwise, to obtain

$$t = R_0^2 \left( \frac{2\rho L}{\gamma a^2} \right)^{1/2} \left[ \frac{\pi}{4} - \frac{\theta}{2} + \frac{1}{4} \sin 2\theta \right],$$

where  $\theta = \sin^{-1}(R/R_0)$  and  $R_0$  is the initial radius of the balloon. Hence find the time when  $R = 0$ .

4/I/7C     **Fluid Dynamics**

Inviscid fluid issues vertically downwards at speed  $u_0$  from a circular tube of radius  $a$ . The fluid falls onto a horizontal plate a distance  $H$  below the end of the tube, where it spreads out axisymmetrically.

Show that while the fluid is falling freely it has speed

$$u = u_0 \left[ 1 + \frac{2g}{u_0^2}(H - z) \right]^{1/2},$$

and occupies a circular jet of radius

$$R = a \left[ 1 + \frac{2g}{u_0^2}(H - z) \right]^{-1/4},$$

where  $z$  is the height above the plate and  $g$  is the acceleration due to gravity.

Show further that along the plate, at radial distances  $r \gg a$  (i.e. far from the falling jet), where the fluid is flowing almost horizontally, it does so as a film of height  $h(r)$ , where

$$\frac{a^4}{4r^2h^2} = 1 + \frac{2g}{u_0^2}(H - h).$$

4/II/16C **Fluid Dynamics**

Define the terms *irrotational flow* and *incompressible flow*. The two-dimensional flow of an incompressible fluid is given in terms of a streamfunction  $\psi(x, y)$  as

$$\mathbf{u} = (u, v) = \left( \frac{\partial \psi}{\partial y}, -\frac{\partial \psi}{\partial x} \right)$$

in Cartesian coordinates  $(x, y)$ . Show that the line integral

$$\int_{\mathbf{x}_1}^{\mathbf{x}_2} \mathbf{u} \cdot \mathbf{n} \, dl = \psi(\mathbf{x}_2) - \psi(\mathbf{x}_1)$$

along any path joining the points  $\mathbf{x}_1$  and  $\mathbf{x}_2$ , where  $\mathbf{n}$  is the unit normal to the path. Describe how this result is related to the concept of mass conservation.

Inviscid, incompressible fluid is contained in the semi-infinite channel  $x > 0$ ,  $0 < y < 1$ , which has rigid walls at  $x = 0$  and at  $y = 0, 1$ , apart from a small opening at the origin through which the fluid is withdrawn with volume flux  $m$  per unit distance in the third dimension. Show that the streamfunction for irrotational flow in the channel can be chosen (up to an additive constant) to satisfy the equation

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

and boundary conditions

$$\begin{aligned} \psi &= 0 && \text{on } y = 0, x > 0, \\ \psi &= -m && \text{on } x = 0, 0 < y < 1, \\ \psi &= -m && \text{on } y = 1, x > 0, \\ \psi &\rightarrow -my && \text{as } x \rightarrow \infty, \end{aligned}$$

if it is assumed that the flow at infinity is uniform. Solve the boundary-value problem above using separation of variables to obtain

$$\psi = -my + \frac{2m}{\pi} \sum_{n=1}^{\infty} \frac{1}{n} \sin n\pi y e^{-n\pi x}.$$

**2/I/4E Further Analysis**

Let  $\tau_1$  be the collection of all subsets  $A \subset \mathbb{N}$  such that  $A = \emptyset$  or  $\mathbb{N} \setminus A$  is finite. Let  $\tau_2$  be the collection of all subsets of  $\mathbb{N}$  of the form  $I_n = \{n, n+1, n+2, \dots\}$ , together with the empty set. Prove that  $\tau_1$  and  $\tau_2$  are both topologies on  $\mathbb{N}$ .

Show that a function  $f$  from the topological space  $(\mathbb{N}, \tau_1)$  to the topological space  $(\mathbb{N}, \tau_2)$  is continuous if and only if one of the following alternatives holds:

- (i)  $f(n) \rightarrow \infty$  as  $n \rightarrow \infty$ ;
- (ii) there exists  $N \in \mathbb{N}$  such that  $f(n) = N$  for all but finitely many  $n$  and  $f(n) \leq N$  for all  $n$ .

**2/II/13E Further Analysis**

(a) Let  $f: [1, \infty) \rightarrow \mathbb{C}$  be defined by  $f(t) = t^{-1}e^{2\pi it}$  and let  $X$  be the image of  $f$ . Prove that  $X \cup \{0\}$  is compact and path-connected. [*Hint: you may find it helpful to set  $s = t^{-1}$ .*]

(b) Let  $g: [1, \infty) \rightarrow \mathbb{C}$  be defined by  $g(t) = (1 + t^{-1})e^{2\pi it}$ , let  $Y$  be the image of  $g$  and let  $\overline{D}$  be the closed unit disc  $\{z \in \mathbb{C} : |z| \leq 1\}$ . Prove that  $Y \cup \overline{D}$  is connected. Explain briefly why it is not path-connected.

**3/I/3E Further Analysis**

(a) Let  $f: \mathbb{C} \rightarrow \mathbb{C}$  be an analytic function such that  $|f(z)| \leq 1 + |z|^{1/2}$  for every  $z$ . Prove that  $f$  is constant.

(b) Let  $f: \mathbb{C} \rightarrow \mathbb{C}$  be an analytic function such that  $\operatorname{Re}(f(z)) \geq 0$  for every  $z$ . Prove that  $f$  is constant.

**3/II/13E Further Analysis**

(a) State Taylor's Theorem.

(b) Let  $f(z) = \sum_{n=0}^{\infty} a_n(z-z_0)^n$  and  $g(z) = \sum_{n=0}^{\infty} b_n(z-z_0)^n$  be defined whenever  $|z-z_0| < r$ . Suppose that  $z_k \rightarrow z_0$  as  $k \rightarrow \infty$ , that no  $z_k$  equals  $z_0$  and that  $f(z_k) = g(z_k)$  for every  $k$ . Prove that  $a_n = b_n$  for every  $n \geq 0$ .

(c) Let  $D$  be a domain, let  $z_0 \in D$  and let  $(z_k)$  be a sequence of points in  $D$  that converges to  $z_0$ , but such that no  $z_k$  equals  $z_0$ . Let  $f: D \rightarrow \mathbb{C}$  and  $g: D \rightarrow \mathbb{C}$  be analytic functions such that  $f(z_k) = g(z_k)$  for every  $k$ . Prove that  $f(z) = g(z)$  for every  $z \in D$ .

(d) Let  $D$  be the domain  $\mathbb{C} \setminus \{0\}$ . Give an example of an analytic function  $f: D \rightarrow \mathbb{C}$  such that  $f(n^{-1}) = 0$  for every positive integer  $n$  but  $f$  is not identically 0.

(e) Show that any function with the property described in (d) must have an essential singularity at the origin.

## 4/I/4E Further Analysis

(a) State and prove Morera's Theorem.

(b) Let  $D$  be a domain and for each  $n \in \mathbb{N}$  let  $f_n : D \rightarrow \mathbb{C}$  be an analytic function. Suppose that  $f : D \rightarrow \mathbb{C}$  is another function and that  $f_n \rightarrow f$  uniformly on  $D$ . Prove that  $f$  is analytic.

## 4/II/13E Further Analysis

(a) State the residue theorem and use it to deduce the principle of the argument, in a form that involves winding numbers.

(b) Let  $p(z) = z^5 + z$ . Find all  $z$  such that  $|z| = 1$  and  $\operatorname{Im}(p(z)) = 0$ . Calculate  $\operatorname{Re}(p(z))$  for each such  $z$ . [It will be helpful to set  $z = e^{i\theta}$ . You may use the addition formulae  $\sin \alpha + \sin \beta = 2 \sin(\frac{\alpha+\beta}{2}) \cos(\frac{\alpha-\beta}{2})$  and  $\cos \alpha + \cos \beta = 2 \cos(\frac{\alpha+\beta}{2}) \cos(\frac{\alpha-\beta}{2})$ .]

(c) Let  $\gamma : [0, 2\pi] \rightarrow \mathbb{C}$  be the closed path  $\theta \mapsto e^{i\theta}$ . Use your answer to (b) to give a rough sketch of the path  $p \circ \gamma$ , paying particular attention to where it crosses the real axis.

(d) Hence, or otherwise, determine for every real  $t$  the number of  $z$  (counted with multiplicity) such that  $|z| < 1$  and  $p(z) = t$ . (You need not give rigorous justifications for your calculations.)

**1/I/4F Geometry**

Describe the geodesics (that is, hyperbolic straight lines) in **either** the disc model **or** the half-plane model of the hyperbolic plane. Explain what is meant by the statements that two hyperbolic lines are parallel, and that they are ultraparallel.

Show that two hyperbolic lines  $l$  and  $l'$  have a unique common perpendicular if and only if they are ultraparallel.

[You may assume standard results about the group of isometries of whichever model of the hyperbolic plane you use.]

**1/II/13F Geometry**

Write down the Riemannian metric in the half-plane model of the hyperbolic plane. Show that Möbius transformations mapping the upper half-plane to itself are isometries of this model.

Calculate the hyperbolic distance from  $ib$  to  $ic$ , where  $b$  and  $c$  are positive real numbers. Assuming that the hyperbolic circle with centre  $ib$  and radius  $r$  is a Euclidean circle, find its Euclidean centre and radius.

Suppose that  $a$  and  $b$  are positive real numbers for which the points  $ib$  and  $a + ib$  of the upper half-plane are such that the hyperbolic distance between them coincides with the Euclidean distance. Obtain an expression for  $b$  as a function of  $a$ . Hence show that, for any  $b$  with  $0 < b < 1$ , there is a unique positive value of  $a$  such that the hyperbolic distance between  $ib$  and  $a + ib$  coincides with the Euclidean distance.

**3/I/4F Geometry**

Show that any isometry of Euclidean 3-space which fixes the origin can be written as a composite of at most three reflections in planes through the origin, and give (with justification) an example of an isometry for which three reflections are necessary.

**3/II/14F Geometry**

State and prove the Gauss–Bonnet formula for the area of a spherical triangle. Deduce a formula for the area of a spherical  $n$ -gon with angles  $\alpha_1, \alpha_2, \dots, \alpha_n$ . For what range of values of  $\alpha$  does there exist a (convex) regular spherical  $n$ -gon with angle  $\alpha$ ?

Let  $\Delta$  be a spherical triangle with angles  $\pi/p, \pi/q$  and  $\pi/r$  where  $p, q, r$  are integers, and let  $G$  be the group of isometries of the sphere generated by reflections in the three sides of  $\Delta$ . List the possible values of  $(p, q, r)$ , and in each case calculate the order of the corresponding group  $G$ . If  $(p, q, r) = (2, 3, 5)$ , show how to construct a regular dodecahedron whose group of symmetries is  $G$ .

[You may assume that the images of  $\Delta$  under the elements of  $G$  form a tessellation of the sphere.]

**1/I/5E Linear Mathematics**

Let  $V$  be the subset of  $\mathbb{R}^5$  consisting of all quintuples  $(a_1, a_2, a_3, a_4, a_5)$  such that

$$a_1 + a_2 + a_3 + a_4 + a_5 = 0$$

and

$$a_1 + 2a_2 + 3a_3 + 4a_4 + 5a_5 = 0 .$$

Prove that  $V$  is a subspace of  $\mathbb{R}^5$ . Solve the above equations for  $a_1$  and  $a_2$  in terms of  $a_3, a_4$  and  $a_5$ . Hence, exhibit a basis for  $V$ , explaining carefully why the vectors you give form a basis.

**1/II/14E Linear Mathematics**

(a) Let  $U, U'$  be subspaces of a finite-dimensional vector space  $V$ . Prove that  $\dim(U + U') = \dim U + \dim U' - \dim(U \cap U')$ .

(b) Let  $V$  and  $W$  be finite-dimensional vector spaces and let  $\alpha$  and  $\beta$  be linear maps from  $V$  to  $W$ . Prove that

$$\text{rank}(\alpha + \beta) \leq \text{rank } \alpha + \text{rank } \beta .$$

(c) Deduce from this result that

$$\text{rank}(\alpha + \beta) \geq |\text{rank } \alpha - \text{rank } \beta| .$$

(d) Let  $V = W = \mathbb{R}^n$  and suppose that  $1 \leq r \leq s \leq n$ . Exhibit linear maps  $\alpha, \beta: V \rightarrow W$  such that  $\text{rank } \alpha = r$ ,  $\text{rank } \beta = s$  and  $\text{rank}(\alpha + \beta) = s - r$ . Suppose that  $r + s \geq n$ . Exhibit linear maps  $\alpha, \beta: V \rightarrow W$  such that  $\text{rank } \alpha = r$ ,  $\text{rank } \beta = s$  and  $\text{rank}(\alpha + \beta) = n$ .

**2/I/6E Linear Mathematics**

Let  $a_1, a_2, \dots, a_n$  be distinct real numbers. For each  $i$  let  $\mathbf{v}_i$  be the vector  $(1, a_i, a_i^2, \dots, a_i^{n-1})$ . Let  $A$  be the  $n \times n$  matrix with rows  $\mathbf{v}_1, \mathbf{v}_2, \dots, \mathbf{v}_n$  and let  $\mathbf{c}$  be a column vector of size  $n$ . Prove that  $A\mathbf{c} = \mathbf{0}$  if and only if  $\mathbf{c} = \mathbf{0}$ . Deduce that the vectors  $\mathbf{v}_1, \mathbf{v}_2, \dots, \mathbf{v}_n$  span  $\mathbb{R}^n$ .

[You may use general facts about matrices if you state them clearly.]

**2/II/15E Linear Mathematics**

(a) Let  $A = (a_{ij})$  be an  $m \times n$  matrix and for each  $k \leq n$  let  $A_k$  be the  $m \times k$  matrix formed by the first  $k$  columns of  $A$ . Suppose that  $n > m$ . Explain why the nullity of  $A$  is non-zero. Prove that if  $k$  is minimal such that  $A_k$  has non-zero nullity, then the nullity of  $A_k$  is 1.

(b) Suppose that no column of  $A$  consists entirely of zeros. Deduce from (a) that there exist scalars  $b_1, \dots, b_k$  (where  $k$  is defined as in (a)) such that  $\sum_{j=1}^k a_{ij}b_j = 0$  for every  $i \leq m$ , but whenever  $\lambda_1, \dots, \lambda_k$  are distinct real numbers there is some  $i \leq m$  such that  $\sum_{j=1}^k a_{ij}\lambda_j b_j \neq 0$ .

(c) Now let  $\mathbf{v}_1, \mathbf{v}_2, \dots, \mathbf{v}_m$  and  $\mathbf{w}_1, \mathbf{w}_2, \dots, \mathbf{w}_m$  be bases for the same real  $m$ -dimensional vector space. Let  $\lambda_1, \lambda_2, \dots, \lambda_n$  be distinct real numbers such that for every  $j$  the vectors  $\mathbf{v}_1 + \lambda_j \mathbf{w}_1, \dots, \mathbf{v}_m + \lambda_j \mathbf{w}_m$  are linearly dependent. For each  $j$ , let  $a_{1j}, \dots, a_{mj}$  be scalars, not all zero, such that  $\sum_{i=1}^m a_{ij}(\mathbf{v}_i + \lambda_j \mathbf{w}_i) = \mathbf{0}$ . By applying the result of (b) to the matrix  $(a_{ij})$ , deduce that  $n \leq m$ .

(d) It follows that the vectors  $\mathbf{v}_1 + \lambda \mathbf{w}_1, \dots, \mathbf{v}_m + \lambda \mathbf{w}_m$  are linearly dependent for at most  $m$  values of  $\lambda$ . Explain briefly how this result can also be proved using determinants.

**3/I/7G Linear Mathematics**

Let  $\alpha$  be an endomorphism of a finite-dimensional real vector space  $U$  and let  $\beta$  be another endomorphism of  $U$  that commutes with  $\alpha$ . If  $\lambda$  is an eigenvalue of  $\alpha$ , show that  $\beta$  maps the kernel of  $\alpha - \lambda \iota$  into itself, where  $\iota$  is the identity map. Suppose now that  $\alpha$  is diagonalizable with  $n$  distinct real eigenvalues where  $n = \dim U$ . Prove that if there exists an endomorphism  $\beta$  of  $U$  such that  $\alpha = \beta^2$ , then  $\lambda \geq 0$  for all eigenvalues  $\lambda$  of  $\alpha$ .

**3/II/17G Linear Mathematics**

Define the *determinant*  $\det(A)$  of an  $n \times n$  complex matrix  $A$ . Let  $A_1, \dots, A_n$  be the columns of  $A$ , let  $\sigma$  be a permutation of  $\{1, \dots, n\}$  and let  $A^\sigma$  be the matrix whose columns are  $A_{\sigma(1)}, \dots, A_{\sigma(n)}$ . Prove from your definition of determinant that  $\det(A^\sigma) = \epsilon(\sigma) \det(A)$ , where  $\epsilon(\sigma)$  is the sign of the permutation  $\sigma$ . Prove also that  $\det(A) = \det(A^t)$ .

Define the *adjugate* matrix  $\text{adj}(A)$  and prove from your definitions that  $A \text{adj}(A) = \text{adj}(A) A = \det(A) I$ , where  $I$  is the identity matrix. Hence or otherwise, prove that if  $\det(A) \neq 0$ , then  $A$  is invertible.

Let  $C$  and  $D$  be real  $n \times n$  matrices such that the complex matrix  $C + iD$  is invertible. By considering  $\det(C + \lambda D)$  as a function of  $\lambda$  or otherwise, prove that there exists a real number  $\lambda$  such that  $C + \lambda D$  is invertible. [You may assume that if a matrix  $A$  is invertible, then  $\det(A) \neq 0$ .]

Deduce that if two real matrices  $A$  and  $B$  are such that there exists an invertible complex matrix  $P$  with  $P^{-1} A P = B$ , then there exists an invertible **real** matrix  $Q$  such that  $Q^{-1} A Q = B$ .



**4/I/6G Linear Mathematics**

Let  $\alpha$  be an endomorphism of a finite-dimensional real vector space  $U$  such that  $\alpha^2 = \alpha$ . Show that  $U$  can be written as the direct sum of the kernel of  $\alpha$  and the image of  $\alpha$ . Hence or otherwise, find the characteristic polynomial of  $\alpha$  in terms of the dimension of  $U$  and the rank of  $\alpha$ . Is  $\alpha$  diagonalizable? Justify your answer.

**4/II/15G Linear Mathematics**

Let  $\alpha \in L(U, V)$  be a linear map between finite-dimensional vector spaces. Let

$$M^l(\alpha) = \{\beta \in L(V, U) : \beta\alpha = 0\} \quad \text{and}$$

$$M^r(\alpha) = \{\beta \in L(V, U) : \alpha\beta = 0\} .$$

(a) Prove that  $M^l(\alpha)$  and  $M^r(\alpha)$  are subspaces of  $L(V, U)$  of dimensions

$$\dim M^l(\alpha) = (\dim V - \text{rank } \alpha) \dim U \quad \text{and}$$

$$\dim M^r(\alpha) = \dim \ker(\alpha) \dim V .$$

[You may use the result that there exist bases in  $U$  and  $V$  so that  $\alpha$  is represented by

$$\begin{pmatrix} I_r & 0 \\ 0 & 0 \end{pmatrix},$$

where  $I_r$  is the  $r \times r$  identity matrix and  $r$  is the rank of  $\alpha$ .]

(b) Let  $\Phi: L(U, V) \rightarrow L(V^*, U^*)$  be given by  $\Phi(\alpha) = \alpha^*$ , where  $\alpha^*$  is the dual map induced by  $\alpha$ . Prove that  $\Phi$  is an isomorphism. [You may assume that  $\Phi$  is linear, and you may use the result that a finite-dimensional vector space and its dual have the same dimension.]

(c) Prove that

$$\Phi(M^l(\alpha)) = M^r(\alpha^*) \quad \text{and} \quad \Phi(M^r(\alpha)) = M^l(\alpha^*).$$

[You may use the results that  $(\beta\alpha)^* = \alpha^*\beta^*$  and that  $\beta^{**}$  can be identified with  $\beta$  under the canonical isomorphism between a vector space and its double dual.]

(d) Conclude that  $\text{rank}(\alpha) = \text{rank}(\alpha^*)$ .

1/I/2D    **Methods**

Fermat's principle of optics states that the path of a light ray connecting two points will be such that the travel time  $t$  is a minimum. If the speed of light varies continuously in a medium and is a function  $c(y)$  of the distance from the boundary  $y = 0$ , show that the path of a light ray is given by the solution to

$$c(y)y'' + c'(y)(1 + y'^2) = 0,$$

where  $y' = \frac{dy}{dx}$ , etc. Show that the path of a light ray in a medium where the speed of light  $c$  is a constant is a straight line. Also find the path from  $(0, 0)$  to  $(1, 0)$  if  $c(y) = y$ , and sketch it.

1/II/11D    **Methods**

(a) Determine the Green's function  $G(x, \xi)$  for the operator  $\frac{d^2}{dx^2} + k^2$  on  $[0, \pi]$  with Dirichlet boundary conditions by solving the boundary value problem

$$\frac{d^2G}{dx^2} + k^2G = \delta(x - \xi), \quad G(0) = 0, \quad G(\pi) = 0$$

when  $k$  is not an integer.

(b) Use the method of Green's functions to solve the boundary value problem

$$\frac{d^2y}{dx^2} + k^2y = f(x), \quad y(0) = a, \quad y(\pi) = b$$

when  $k$  is not an integer.

**2/I/2C Methods**

Explain briefly why the second-rank tensor

$$\int_S x_i x_j dS(\mathbf{x})$$

is isotropic, where  $S$  is the surface of the unit sphere centred on the origin.

A second-rank tensor is defined by

$$T_{ij}(\mathbf{y}) = \int_S (y_i - x_i)(y_j - x_j) dS(\mathbf{x}),$$

where  $S$  is the surface of the unit sphere centred on the origin. Calculate  $T(\mathbf{y})$  in the form

$$T_{ij} = \lambda \delta_{ij} + \mu y_i y_j,$$

where  $\lambda$  and  $\mu$  are to be determined.

By considering the action of  $T$  on  $\mathbf{y}$  and on vectors perpendicular to  $\mathbf{y}$ , determine the eigenvalues and associated eigenvectors of  $T$ .

**2/II/11C Methods**

State the transformation law for an  $n$ th-rank tensor  $T_{ij\dots k}$ .

Show that the fourth-rank tensor

$$c_{ijkl} = \alpha \delta_{ij} \delta_{kl} + \beta \delta_{ik} \delta_{jl} + \gamma \delta_{il} \delta_{jk}$$

is isotropic for arbitrary scalars  $\alpha$ ,  $\beta$  and  $\gamma$ .

The stress  $\sigma_{ij}$  and strain  $e_{ij}$  in a linear elastic medium are related by

$$\sigma_{ij} = c_{ijkl} e_{kl}.$$

Given that  $e_{ij}$  is symmetric and that the medium is isotropic, show that the stress-strain relationship can be written in the form

$$\sigma_{ij} = \lambda e_{kk} \delta_{ij} + 2\mu e_{ij}.$$

Show that  $e_{ij}$  can be written in the form  $e_{ij} = p\delta_{ij} + d_{ij}$ , where  $d_{ij}$  is a traceless tensor and  $p$  is a scalar to be determined. Show also that necessary and sufficient conditions for the stored elastic energy density  $E = \frac{1}{2}\sigma_{ij} e_{ij}$  to be non-negative for any deformation of the solid are that

$$\mu \geq 0 \quad \text{and} \quad \lambda \geq -\frac{2}{3}\mu.$$

### 3/I/2D Methods

Consider the path between two arbitrary points on a cone of interior angle  $2\alpha$ . Show that the arc-length of the path  $r(\theta)$  is given by

$$\int (r^2 + r'^2 \operatorname{cosec}^2 \alpha)^{1/2} d\theta,$$

where  $r' = \frac{dr}{d\theta}$ . By minimizing the total arc-length between the points, determine the equation for the shortest path connecting them.

### 3/II/12D Methods

The transverse displacement  $y(x, t)$  of a stretched string clamped at its ends  $x = 0, l$  satisfies the equation

$$\frac{\partial^2 y}{\partial t^2} = c^2 \frac{\partial^2 y}{\partial x^2} - 2k \frac{\partial y}{\partial t}, \quad y(x, 0) = 0, \quad \frac{\partial y}{\partial t}(x, 0) = \delta(x - a),$$

where  $c > 0$  is the wave velocity, and  $k > 0$  is the damping coefficient. The initial conditions correspond to a sharp blow at  $x = a$  at time  $t = 0$ .

(a) Show that the subsequent motion of the string is given by

$$y(x, t) = \frac{1}{\sqrt{\alpha_n^2 - k^2}} \sum_n 2e^{-kt} \sin \frac{\alpha_n a}{c} \sin \frac{\alpha_n x}{c} \sin /(\sqrt{\alpha_n^2 - k^2} t)$$

where  $\alpha_n = \pi cn/l$ .

(b) Describe what happens in the limits of small and large damping. What critical parameter separates the two cases?

### 4/I/2D Methods

Consider the wave equation in a spherically symmetric coordinate system

$$\frac{\partial^2 u(r, t)}{\partial t^2} = c^2 \Delta u(r, t),$$

where  $\Delta u = \frac{1}{r} \frac{\partial^2}{\partial r^2} (ru)$  is the spherically symmetric Laplacian operator.

(a) Show that the general solution to the equation above is

$$u(r, t) = \frac{1}{r} [f(r + ct) + g(r - ct)],$$

where  $f(x), g(x)$  are arbitrary functions.

(b) Using separation of variables, determine the wave field  $u(r, t)$  in response to a pulsating source at the origin  $u(0, t) = A \sin \omega t$ .

#### 4/II/11D Methods

The velocity potential  $\phi(r, \theta)$  for inviscid flow in two dimensions satisfies the Laplace equation

$$\Delta\phi = \left[ \frac{1}{r} \frac{\partial}{\partial r} \left( r \frac{\partial}{\partial r} \right) + \frac{1}{r^2} \frac{\partial^2}{\partial \theta^2} \right] \phi(r, \theta) = 0 .$$

(a) Using separation of variables, derive the general solution to the equation above that is single-valued and finite in each of the domains (i)  $0 \leq r \leq a$ ; (ii)  $a \leq r < \infty$ .

(b) Assuming  $\phi$  is single-valued, solve the Laplace equation subject to the boundary conditions  $\frac{\partial\phi}{\partial r} = 0$  at  $r = a$ , and  $\frac{\partial\phi}{\partial r} \rightarrow U \cos\theta$  as  $r \rightarrow \infty$ . Sketch the lines of constant potential.

**2/I/5B Numerical Analysis**

Let

$$A = \begin{pmatrix} 1 & a & a^2 & a^3 \\ a^3 & 1 & a & a^2 \\ a^2 & a^3 & 1 & a \\ a & a^2 & a^3 & 1 \end{pmatrix}, \quad b = \begin{pmatrix} \gamma \\ 0 \\ 0 \\ \gamma a \end{pmatrix}, \quad \gamma = 1 - a^4 \neq 0.$$

Find the LU factorization of the matrix  $A$  and use it to solve the system  $Ax = b$ .

**2/II/14B Numerical Analysis**

Let

$$f''(0) \approx a_0 f(-1) + a_1 f(0) + a_2 f(1) = \mu(f)$$

be an approximation of the second derivative which is exact for  $f \in \mathcal{P}_2$ , the set of polynomials of degree  $\leq 2$ , and let

$$e(f) = f''(0) - \mu(f)$$

be its error.

(a) Determine the coefficients  $a_0, a_1, a_2$ .

(b) Using the Peano kernel theorem prove that, for  $f \in C^3[-1, 1]$ , the set of three-times continuously differentiable functions, the error satisfies the inequality

$$|e(f)| \leq \frac{1}{3} \max_{x \in [-1, 1]} |f'''(x)|.$$

**3/I/6B Numerical Analysis**

Given  $(n + 1)$  distinct points  $x_0, x_1, \dots, x_n$ , let

$$\ell_i(x) = \prod_{\substack{k=0 \\ k \neq i}}^n \frac{x - x_k}{x_i - x_k}$$

be the fundamental Lagrange polynomials of degree  $n$ , let

$$\omega(x) = \prod_{i=0}^n (x - x_i),$$

and let  $p$  be any polynomial of degree  $\leq n$ .

(a) Prove that  $\sum_{i=0}^n p(x_i)\ell_i(x) \equiv p(x)$ .

(b) Hence or otherwise derive the formula

$$\frac{p(x)}{\omega(x)} = \sum_{i=0}^n \frac{A_i}{x - x_i}, \quad A_i = \frac{p(x_i)}{\omega'(x_i)},$$

which is the decomposition of  $p(x)/\omega(x)$  into partial fractions.

**3/II/16B Numerical Analysis**

The functions  $H_0, H_1, \dots$  are generated by the Rodrigues formula:

$$H_n(x) = (-1)^n e^{x^2} \frac{d^n}{dx^n} e^{-x^2}.$$

(a) Show that  $H_n$  is a polynomial of degree  $n$ , and that the  $H_n$  are orthogonal with respect to the scalar product

$$(f, g) = \int_{-\infty}^{\infty} f(x)g(x)e^{-x^2} dx.$$

(b) By induction or otherwise, prove that the  $H_n$  satisfy the three-term recurrence relation

$$H_{n+1}(x) = 2xH_n(x) - 2nH_{n-1}(x).$$

[Hint: you may need to prove the equality  $H_n'(x) = 2nH_{n-1}(x)$  as well.]

3/I/5H **Optimization**

Two players A and B play a zero-sum game with the pay-off matrix

	$B_1$	$B_2$	$B_3$
$A_1$	4	-2	-5
$A_2$	-2	4	3
$A_3$	-3	6	2
$A_4$	3	-8	-6

Here, the  $(i, j)$  entry of the matrix indicates the pay-off to player A if he chooses move  $A_i$  and player B chooses move  $B_j$ . Show that the game can be reduced to a zero-sum game with  $2 \times 2$  pay-off matrix.

Determine the value of the game and the optimal strategy for player A.

3/II/15H **Optimization**

Explain what is meant by a transportation problem where the total demand equals the total supply. Write the Lagrangian and describe an algorithm for solving such a problem. Starting from the north-west initial assignment, solve the problem with three sources and three destinations described by the table

5	9	1	36
3	10	6	84
7	2	5	40
14	68	78	

where the figures in the  $3 \times 3$  box denote the transportation costs (per unit), the right-hand column denotes supplies, and the bottom row demands.

4/I/5H **Optimization**

State and prove the Lagrangian sufficiency theorem for a general optimization problem with constraints.

4/II/14H **Optimization**

Use the two-phase simplex method to solve the problem

$$\begin{aligned}
 &\text{minimize} && 5x_1 - 12x_2 + 13x_3 \\
 &\text{subject to} && 4x_1 + 5x_2 \leq 9, \\
 &&& 6x_1 + 4x_2 + x_3 \geq 12, \\
 &&& 3x_1 + 2x_2 - x_3 \leq 3, \\
 &&& x_i \geq 0, \quad i = 1, 2, 3.
 \end{aligned}$$



**1/I/8G Quadratic Mathematics**

Let  $U$  and  $V$  be finite-dimensional vector spaces. Suppose that  $b$  and  $c$  are bilinear forms on  $U \times V$  and that  $b$  is non-degenerate. Show that there exist linear endomorphisms  $S$  of  $U$  and  $T$  of  $V$  such that  $c(x, y) = b(S(x), y) = b(x, T(y))$  for all  $(x, y) \in U \times V$ .

**1/II/17G Quadratic Mathematics**

(a) Suppose  $p$  is an odd prime and  $a$  an integer coprime to  $p$ . Define the *Legendre symbol*  $\left(\frac{a}{p}\right)$  and state Euler's criterion.

(b) Compute  $\left(\frac{-1}{p}\right)$  and prove that

$$\left(\frac{ab}{p}\right) = \left(\frac{a}{p}\right) \left(\frac{b}{p}\right)$$

whenever  $a$  and  $b$  are coprime to  $p$ .

(c) Let  $n$  be any integer such that  $1 \leq n \leq p-2$ . Let  $m$  be the unique integer such that  $1 \leq m \leq p-2$  and  $mn \equiv 1 \pmod{p}$ . Prove that

$$\left(\frac{n(n+1)}{p}\right) = \left(\frac{1+m}{p}\right).$$

(d) Find

$$\sum_{n=1}^{p-2} \left(\frac{n(n+1)}{p}\right).$$

**2/I/8G Quadratic Mathematics**

Let  $U$  be a finite-dimensional real vector space and  $b$  a positive definite symmetric bilinear form on  $U \times U$ . Let  $\psi: U \rightarrow U$  be a linear map such that  $b(\psi(x), y) + b(x, \psi(y)) = 0$  for all  $x$  and  $y$  in  $U$ . Prove that if  $\psi$  is invertible, then the dimension of  $U$  must be even. By considering the restriction of  $\psi$  to its image or otherwise, prove that the rank of  $\psi$  is always even.

**2/II/17G Quadratic Mathematics**

Let  $S$  be the set of all  $2 \times 2$  complex matrices  $A$  which are *hermitian*, that is,  $A^* = A$ , where  $A^* = \overline{A}^t$ .

(a) Show that  $S$  is a real 4-dimensional vector space. Consider the real symmetric bilinear form  $b$  on this space defined by

$$b(A, B) = \frac{1}{2} (\operatorname{tr}(AB) - \operatorname{tr}(A) \operatorname{tr}(B)) .$$

Prove that  $b(A, A) = -\det A$  and  $b(A, I) = -\frac{1}{2}\operatorname{tr}(A)$ , where  $I$  denotes the identity matrix.

(b) Consider the three matrices

$$A_1 = \begin{pmatrix} 1 & 0 \\ 0 & -1 \end{pmatrix} , \quad A_2 = \begin{pmatrix} 0 & 1 \\ 1 & 0 \end{pmatrix} \quad \text{and} \quad A_3 = \begin{pmatrix} 0 & -i \\ i & 0 \end{pmatrix} .$$

Prove that the basis  $I, A_1, A_2, A_3$  of  $S$  diagonalizes  $b$ . Hence or otherwise find the rank and signature of  $b$ .

(c) Let  $Q$  be the set of all  $2 \times 2$  complex matrices  $C$  which satisfy  $C + C^* = \operatorname{tr}(C) I$ . Show that  $Q$  is a real 4-dimensional vector space. Given  $C \in Q$ , put

$$\Phi(C) = \frac{1-i}{2} \operatorname{tr}(C) I + i C .$$

Show that  $\Phi$  takes values in  $S$  and is a linear isomorphism between  $Q$  and  $S$ .

(d) Define a real symmetric bilinear form on  $Q$  by setting  $c(C, D) = -\frac{1}{2}\operatorname{tr}(CD)$ ,  $C, D \in Q$ . Show that  $b(\Phi(C), \Phi(D)) = c(C, D)$  for all  $C, D \in Q$ . Find the rank and signature of the symmetric bilinear form  $c$  defined on  $Q$ .

**3/I/9G Quadratic Mathematics**

Let  $f(x, y) = ax^2 + bxy + cy^2$  be a binary quadratic form with integer coefficients. Explain what is meant by the *discriminant*  $d$  of  $f$ . State a necessary and sufficient condition for some form of discriminant  $d$  to represent an odd prime number  $p$ . Using this result or otherwise, find the primes  $p$  which can be represented by the form  $x^2 + 3y^2$ .

3/II/19G **Quadratic Mathematics**

Let  $U$  be a finite-dimensional real vector space endowed with a positive definite inner product. A linear map  $\tau : U \rightarrow U$  is said to be an *orthogonal projection* if  $\tau$  is self-adjoint and  $\tau^2 = \tau$ .

(a) Prove that for every orthogonal projection  $\tau$  there is an orthogonal decomposition

$$U = \ker(\tau) \oplus \operatorname{im}(\tau).$$

(b) Let  $\phi : U \rightarrow U$  be a linear map. Show that if  $\phi^2 = \phi$  and  $\phi\phi^* = \phi^*\phi$ , where  $\phi^*$  is the adjoint of  $\phi$ , then  $\phi$  is an orthogonal projection. [*You may find it useful to prove first that if  $\phi\phi^* = \phi^*\phi$ , then  $\phi$  and  $\phi^*$  have the same kernel.*]

(c) Show that given a subspace  $W$  of  $U$  there exists a unique orthogonal projection  $\tau$  such that  $\operatorname{im}(\tau) = W$ . If  $W_1$  and  $W_2$  are two subspaces with corresponding orthogonal projections  $\tau_1$  and  $\tau_2$ , show that  $\tau_2 \circ \tau_1 = 0$  if and only if  $W_1$  is orthogonal to  $W_2$ .

(d) Let  $\phi : U \rightarrow U$  be a linear map satisfying  $\phi^2 = \phi$ . Prove that one can define a positive definite inner product on  $U$  such that  $\phi$  becomes an orthogonal projection.

**1/I/9A Quantum Mechanics**

A particle of mass  $m$  is confined inside a one-dimensional box of length  $a$ . Determine the possible energy eigenvalues.

**1/II/18A Quantum Mechanics**

What is the significance of the expectation value

$$\langle Q \rangle = \int \psi^*(x) Q \psi(x) dx$$

of an observable  $Q$  in the normalized state  $\psi(x)$ ? Let  $Q$  and  $P$  be two observables. By considering the norm of  $(Q + i\lambda P)\psi$  for real values of  $\lambda$ , show that

$$\langle Q^2 \rangle \langle P^2 \rangle \geq \frac{1}{4} |\langle [Q, P] \rangle|^2.$$

The uncertainty  $\Delta Q$  of  $Q$  in the state  $\psi(x)$  is defined as

$$(\Delta Q)^2 = \langle (Q - \langle Q \rangle)^2 \rangle.$$

Deduce the generalized uncertainty relation,

$$\Delta Q \Delta P \geq \frac{1}{2} |\langle [Q, P] \rangle|.$$

A particle of mass  $m$  moves in one dimension under the influence of the potential  $\frac{1}{2}m\omega^2 x^2$ . By considering the commutator  $[x, p]$ , show that the expectation value of the Hamiltonian satisfies

$$\langle H \rangle \geq \frac{1}{2} \hbar \omega.$$

**2/I/9A Quantum Mechanics**

What is meant by the statement that an operator is *hermitian*?

A particle of mass  $m$  moves in the real potential  $V(x)$  in one dimension. Show that the Hamiltonian of the system is hermitian.

Show that

$$\begin{aligned} \frac{d}{dt} \langle x \rangle &= \frac{1}{m} \langle p \rangle, \\ \frac{d}{dt} \langle p \rangle &= \langle -V'(x) \rangle, \end{aligned}$$

where  $p$  is the momentum operator and  $\langle A \rangle$  denotes the expectation value of the operator  $A$ .

2/II/18A **Quantum Mechanics**

A particle of mass  $m$  and energy  $E$  moving in one dimension is incident from the left on a potential barrier  $V(x)$  given by

$$V(x) = \begin{cases} V_0 & 0 \leq x \leq a \\ 0 & \text{otherwise} \end{cases}$$

with  $V_0 > 0$ .

In the limit  $V_0 \rightarrow \infty, a \rightarrow 0$  with  $V_0 a = U$  held fixed, show that the transmission probability is

$$T = \left(1 + \frac{mU^2}{2E\hbar^2}\right)^{-1}.$$

3/II/20A **Quantum Mechanics**

The radial wavefunction for the hydrogen atom satisfies the equation

$$\frac{-\hbar^2}{2m} \frac{1}{r^2} \frac{d}{dr} \left( r^2 \frac{d}{dr} R(r) \right) + \frac{\hbar^2}{2mr^2} \ell(\ell+1) R(r) - \frac{e^2}{4\pi\epsilon_0 r} R(r) = ER(r).$$

Explain the origin of each term in this equation.

The wavefunctions for the ground state and first radially excited state, both with  $\ell = 0$ , can be written as

$$R_1(r) = N_1 \exp(-\alpha r)$$

$$R_2(r) = N_2(r+b) \exp(-\beta r)$$

respectively, where  $N_1$  and  $N_2$  are normalization constants. Determine  $\alpha, \beta, b$  and the corresponding energy eigenvalues  $E_1$  and  $E_2$ .

A hydrogen atom is in the first radially excited state. It makes the transition to the ground state, emitting a photon. What is the frequency of the emitted photon?

**3/I/10A Special Relativity**

What are the momentum and energy of a photon of wavelength  $\lambda$ ?

A photon of wavelength  $\lambda$  is incident on an electron. After the collision, the photon has wavelength  $\lambda'$ . Show that

$$\lambda' - \lambda = \frac{h}{mc}(1 - \cos \theta),$$

where  $\theta$  is the scattering angle and  $m$  is the electron rest mass.

**4/I/9A Special Relativity**

Prove that the two-dimensional Lorentz transformation can be written in the form

$$\begin{aligned}x' &= x \cosh \phi - ct \sinh \phi \\ct' &= -x \sinh \phi + ct \cosh \phi,\end{aligned}$$

where  $\tanh \phi = v/c$ . Hence, show that

$$\begin{aligned}x' + ct' &= e^{-\phi}(x + ct) \\x' - ct' &= e^{\phi}(x - ct).\end{aligned}$$

Given that frame  $S'$  has speed  $v$  with respect to  $S$  and  $S''$  has speed  $v'$  with respect to  $S'$ , use this formalism to find the speed  $v''$  of  $S''$  with respect to  $S$ .

[*Hint: rotation through a hyperbolic angle  $\phi$ , followed by rotation through  $\phi'$ , is equivalent to rotation through  $\phi + \phi'$ .*]

**4/II/18A Special Relativity**

A pion of rest mass  $M$  decays at rest into a muon of rest mass  $m < M$  and a neutrino of zero rest mass. What is the speed  $u$  of the muon?

In the pion rest frame  $S$ , the muon moves in the  $y$ -direction. A moving observer, in a frame  $S'$  with axes parallel to those in the pion rest frame, wishes to take measurements of the decay along the  $x$ -axis, and notes that the pion has speed  $v$  with respect to the  $x$ -axis. Write down the four-dimensional Lorentz transformation relating  $S'$  to  $S$  and determine the momentum of the muon in  $S'$ . Hence show that in  $S'$  the direction of motion of the muon makes an angle  $\theta$  with respect to the  $y$ -axis, where

$$\tan \theta = \frac{M^2 + m^2}{M^2 - m^2} \frac{v}{(c^2 - v^2)^{1/2}}.$$

**1/I/3H Statistics**

Derive the least squares estimators  $\hat{\alpha}$  and  $\hat{\beta}$  for the coefficients of the simple linear regression model

$$Y_i = \alpha + \beta(x_i - \bar{x}) + \varepsilon_i, \quad i = 1, \dots, n,$$

where  $x_1, \dots, x_n$  are given constants,  $\bar{x} = n^{-1} \sum_{i=1}^n x_i$ , and  $\varepsilon_i$  are independent with  $E \varepsilon_i = 0$ ,  $\text{Var } \varepsilon_i = \sigma^2$ ,  $i = 1, \dots, n$ .

A manufacturer of optical equipment has the following data on the unit cost (in pounds) of certain custom-made lenses and the number of units made in each order:

<i>No. of units, <math>x_i</math></i>	1	3	5	10	12
<i>Cost per unit, <math>y_i</math></i>	58	55	40	37	22

Assuming that the conditions underlying simple linear regression analysis are met, estimate the regression coefficients and use the estimated regression equation to predict the unit cost in an order for 8 of these lenses.

[*Hint: for the data above,  $S_{xy} = \sum_{i=1}^n (x_i - \bar{x})y_i = -257.4$ .*]

**1/II/12H Statistics**

Suppose that six observations  $X_1, \dots, X_6$  are selected at random from a normal distribution for which both the mean  $\mu_X$  and the variance  $\sigma_X^2$  are unknown, and it is found that  $S_{XX} = \sum_{i=1}^6 (x_i - \bar{x})^2 = 30$ , where  $\bar{x} = \frac{1}{6} \sum_{i=1}^6 x_i$ . Suppose also that 21 observations  $Y_1, \dots, Y_{21}$  are selected at random from another normal distribution for which both the mean  $\mu_Y$  and the variance  $\sigma_Y^2$  are unknown, and it is found that  $S_{YY} = 40$ . Derive carefully the likelihood ratio test of the hypothesis  $H_0: \sigma_X^2 = \sigma_Y^2$  against  $H_1: \sigma_X^2 > \sigma_Y^2$  and apply it to the data above at the 0.05 level.

[*Hint:*

<i>Distribution</i>	$\chi_5^2$	$\chi_6^2$	$\chi_{20}^2$	$\chi_{21}^2$	$F_{5,20}$	$F_{6,21}$
<i>95% percentile</i>	11.07	12.59	31.41	32.68	2.71	2.57 ]

**2/I/3H Statistics**

Let  $X_1, \dots, X_n$  be a random sample from the  $N(\theta, \sigma^2)$  distribution, and suppose that the prior distribution for  $\theta$  is  $N(\mu, \tau^2)$ , where  $\sigma^2, \mu, \tau^2$  are known. Determine the posterior distribution for  $\theta$ , given  $X_1, \dots, X_n$ , and the best point estimate of  $\theta$  under both quadratic and absolute error loss.

**2/II/12H Statistics**

An examination was given to 500 high-school students in each of two large cities, and their grades were recorded as low, medium, or high. The results are given in the table below.

	<i>Low</i>	<i>Medium</i>	<i>High</i>
<i>City A</i>	103	145	252
<i>City B</i>	140	136	224

Derive carefully the test of homogeneity and test the hypothesis that the distributions of scores among students in the two cities are the same.

[*Hint:*

<i>Distribution</i>	$\chi_1^2$	$\chi_2^2$	$\chi_3^2$	$\chi_5^2$	$\chi_6^2$
<i>99% percentile</i>	6.63	9.21	11.34	15.09	16.81
<i>95% percentile</i>	3.84	5.99	7.81	11.07	12.59

**4/I/3H Statistics**

The following table contains a distribution obtained in 320 tosses of 6 coins and the corresponding expected frequencies calculated with the formula for the binomial distribution for  $p = 0.5$  and  $n = 6$ .

No. heads	0	1	2	3	4	5	6
Observed frequencies	3	21	85	110	62	32	7
Expected frequencies	5	30	75	100	75	30	5

Conduct a goodness-of-fit test at the 0.05 level for the null hypothesis that the coins are all fair.

[*Hint:*

<i>Distribution</i>	$\chi_5^2$	$\chi_6^2$	$\chi_7^2$
<i>95% percentile</i>	11.07	12.59	14.07

**4/II/12H Statistics**

State and prove the Rao–Blackwell theorem.

Suppose that  $X_1, \dots, X_n$  are independent random variables uniformly distributed over  $(\theta, 3\theta)$ . Find a two-dimensional sufficient statistic  $T(X)$  for  $\theta$ . Show that an unbiased estimator of  $\theta$  is  $\hat{\theta} = X_1/2$ .

Find an unbiased estimator of  $\theta$  which is a function of  $T(X)$  and whose mean square error is no more than that of  $\hat{\theta}$ .