# MATHEMATICAL TRIPOS Part II

Friday, 11 September, 2020 1:30 pm to 4:30 pm

# PAPER 4

# Before you begin read these instructions carefully.

The examination paper is divided into two sections. Each question in Section II carries twice the number of marks of each question in Section I. Section II questions also carry an alpha or beta quality mark and Section I questions carry a beta quality mark.

Candidates may obtain credit from attempts on **at most six questions** from Section I and from any number of questions from Section II.

Write on **one side** of the paper only and begin each answer on a separate sheet. Write legibly; otherwise you place yourself at a grave disadvantage.

# At the end of the examination:

Separate your answers to each question.

Complete a gold cover sheet **for each question** that you have attempted, and place it at the front of your answer to that question.

Complete a green master cover sheet listing **all the questions** that you have attempted.

Every cover sheet must also show your Blind Grade ID and desk number.

Tie up your answers and cover sheets into **a single bundle**, with the master cover sheet on the top, and then the cover sheet and answer for each question, in the numerical order of the questions.

# STATIONERY REQUIREMENTS

Gold cover sheets Green master cover sheet Script paper Rough paper

> You may not start to read the questions printed on the subsequent pages until instructed to do so by the Invigilator.

# SECTION I

**1H** Number Theory Let p be a prime.

State and prove Lagrange's theorem on the number of solutions of a polynomial congruence modulo p. Deduce that  $(p-1)! \equiv -1 \mod p$ .

Let k be a positive integer such that k|(p-1). Show that the congruence

 $x^k \equiv 1 \mod p$ 

has precisely k solutions modulo p.

#### 2H Topics in Analysis

Define what is meant by a *nowhere dense* set in a metric space. State a version of the Baire Category theorem.

Let  $f: [1, \infty) \to \mathbb{R}$  be a continuous function such that  $f(nx) \to 0$  as  $n \to \infty$  for every fixed  $x \ge 1$ . Show that  $f(t) \to 0$  as  $t \to \infty$ .

#### 3I Coding and Cryptography

(a) What does it mean to say that a cipher has *perfect secrecy*? Show that if a cipher has perfect secrecy then there must be at least as many possible keys as there are possible plaintext messages. What is a *one-time pad*? Show that a one-time pad has perfect secrecy.

(b) I encrypt a binary sequence  $a_1, a_2, \ldots, a_N$  using a one-time pad with key sequence  $k_1, k_2, k_3, \ldots$  I transmit  $a_1 + k_1, a_2 + k_2, \ldots, a_N + k_N$  to you. Then, by mistake, I also transmit  $a_1 + k_2, a_2 + k_3, \ldots, a_N + k_{N+1}$  to you. Assuming that you know I have made this error, and that my message makes sense, how would you go about finding my message? Can you now decipher other messages sent using the same part of the key sequence? Briefly justify your answer.

#### 4F Automata and Formal Languages

Define what it means for a context-free grammar (CFG) to be in *Chomsky normal form* (CNF).

Describe without proof each stage in the process of converting a CFG  $G = (N, \Sigma, P, S)$  into an equivalent CFG  $\overline{G}$  which is in CNF. For each of these stages, when are the nonterminals N left unchanged? What about the terminals  $\Sigma$  and the generated language  $\mathcal{L}(G)$ ?

Give an example of a CFG G whose generated language  $\mathcal{L}(G)$  is infinite and equal to  $\mathcal{L}(\overline{G})$ .

### 5J Statistical Modelling

Suppose you have a data frame with variables <code>response</code>, <code>covar1</code>, and <code>covar2</code>. You run the following commands on R.

```
model <- lm(response ~ covar1 + covar2)
summary(model)
...
Estimate Std. Error t value Pr(>|t|)
(Intercept) -2.1024 0.1157 -18.164 <2e-16
covar1 1.6329 2.6557 0.615 0.542
covar2 0.3755 2.5978 0.145 0.886</pre>
```

• • •

(a) Consider the following three scenarios:

(i) All the output you have is the abbreviated output of summary(model) above.

(ii) You have the abbreviated output of summary(model) above together with

Residual standard error: 0.8097 on 47 degrees of freedom Multiple R-squared: 0.8126, Adjusted R-squared: 0.8046 F-statistic: 101.9 on 2 and 47 DF, p-value: < 2.2e-16

(iii) You have the abbreviated output of summary(model) above together with

Residual standard error: 0.9184 on 47 degrees of freedom Multiple R-squared: 0.000712, Adjusted R-squared: -0.04181 F-statistic: 0.01674 on 2 and 47 DF, p-value: 0.9834

What conclusion can you draw about which variables explain the response in each of the three scenarios? Explain.

(b) Assume now that you have the abbreviated output of summary(model) above together with

anova(lm(response ~ 1), lm(response ~ covar1), model) . . . F Pr(>F) Res.Df RSS Df Sum of Sq 49 164.448 1 ? <2e-16 2 ? 30.831 ? 133.618 3 ? 30.817 ? 0.014 ? ? . . .

What are the values of the entries with a question mark? [You may express your answers as arithmetic expressions if necessary].

#### 6B Mathematical Biology

Consider a population process in which the probability of transition from a state with n individuals to a state with n + 1 individuals in the interval  $(t, t + \Delta t)$  is  $\lambda n \Delta t$  for small  $\Delta t$ .

(i) Write down the master equation for the probability,  $P_n(t)$ , of the state n at time t for  $n \ge 1$ .

(ii) Assuming an initial distribution

$$P_n(0) = \begin{cases} 1, & \text{if } n = 1, \\ 0, & \text{if } n > 1, \end{cases}$$

show that

$$P_n(t) = \exp(-\lambda t)(1 - \exp(-\lambda t))^{n-1}.$$

(iii) Hence, determine the mean of n for t > 0.

#### 7E Further Complex Methods

The Hilbert transform of a function f(x) is defined by

$$\mathcal{H}(f)(y) := \frac{1}{\pi} \mathcal{P} \int_{-\infty}^{+\infty} \frac{f(x)}{y - x} dx .$$

Calculate the Hilbert transform of  $f(x) = \cos \omega x$ , where  $\omega$  is a non-zero real constant.

#### 8B Classical Dynamics

Derive expressions for the angular momentum and kinetic energy of a rigid body in terms of its mass M, the position  $\mathbf{X}(t)$  of its centre of mass, its inertia tensor I (which should be defined) about its centre of mass, and its angular velocity  $\boldsymbol{\omega}$ .

A spherical planet of mass M and radius R has density proportional to  $r^{-1}\sin(\pi r/R)$ . Given that  $\int_0^{\pi} x \sin x \, dx = \pi$  and  $\int_0^{\pi} x^3 \sin x \, dx = \pi(\pi^2 - 6)$ , evaluate the inertia tensor of the planet in terms of M and R.

#### 9D Cosmology

At temperature T and chemical potential  $\mu$ , the number density of a non-relativistic particle species with mass  $m \gg k_B T/c^2$  is given by

$$n = g \left(\frac{mk_BT}{2\pi\hbar^2}\right)^{3/2} e^{-(mc^2-\mu)/k_BT} \,,$$

where g is the number of degrees of freedom of this particle.

At recombination, electrons and protons combine to form hydrogen. Use the result above to derive the Saha equation

$$n_H \approx n_e^2 \left(\frac{2\pi\hbar^2}{m_e k_B T}\right)^{3/2} e^{E_{\rm bind}/k_B T} \,,$$

where  $n_H$  is the number density of hydrogen atoms,  $n_e$  the number density of electrons,  $m_e$  the mass of the electron and  $E_{\text{bind}}$  the binding energy of hydrogen. State any assumptions that you use in this derivation.

#### 10C Quantum Information and Computation

(i) What is the action of  $QFT_N$  on a state  $|x\rangle$ , where  $x \in \{0, 1, 2, ..., N-1\}$  and  $QFT_N$  denotes the Quantum Fourier Transform modulo N?

(ii) For the case N = 4 write 0, 1, 2, 3 respectively in binary as 00, 01, 10, 11 thereby identifying the four-dimensional space as that of two qubits. Show that  $QFT_N |10\rangle$  is an unentangled state of the two qubits.

(iii) Prove that  $(\operatorname{QFT}_N)^2 |x\rangle = |N - x\rangle$ , where  $(\operatorname{QFT}_N)^2 \equiv \operatorname{QFT}_N \circ \operatorname{QFT}_N$ . [*Hint: For*  $\omega = e^{2\pi i/N}$ ,  $\sum_{m=0}^{N-1} \omega^{mK} = 0$  if K is not a multiple of N.]

(iv) What is the action of  $(QFT_N)^4$  on a state  $|x\rangle$ , for any  $x \in \{0, 1, 2, ..., N-1\}$ ? Use the above to determine what the eigenvalues of  $QFT_N$  are.

# SECTION II

## 11H Number Theory

(a) What does it mean to say that a function  $f : \mathbb{N} \to \mathbb{C}$  is *multiplicative*? Show that if  $f, g : \mathbb{N} \to \mathbb{C}$  are both multiplicative, then so is  $f \star g : \mathbb{N} \to \mathbb{C}$ , defined for all  $n \in \mathbb{N}$  by

$$f \star g(n) = \sum_{d|n} f(d) g\left(\frac{n}{d}\right).$$

Show that if  $f = \mu \star g$ , where  $\mu$  is the Möbius function, then  $g = f \star 1$ , where 1 denotes the constant function 1.

(b) Let  $\tau(n)$  denote the number of positive divisors of n. Find  $f, g: \mathbb{N} \to \mathbb{C}$  such that  $\tau = f \star g$ , and deduce that  $\tau$  is multiplicative. Hence or otherwise show that for all  $s \in \mathbb{C}$  with  $\operatorname{Re}(s) > 1$ ,

$$\sum_{n=1}^{\infty} \frac{\tau(n)}{n^s} = \zeta(s)^2,$$

where  $\zeta$  is the Riemann zeta function.

(c) Fix  $k \in \mathbb{N}$ . By considering suitable powers of the product of the first k + 1 primes, show that

$$\tau(n) \ge (\log n)^k$$

for infinitely many  $n \in \mathbb{N}$ .

(d) Fix  $\epsilon > 0$ . Show that

$$\frac{\tau(n)}{n^{\epsilon}} = \prod_{p \text{ prime, } p^{\alpha} \mid \mid n} \frac{(\alpha + 1)}{p^{\alpha \epsilon}},$$

where  $p^{\alpha} \mid \mid n$  denotes the fact that  $\alpha \in \mathbb{N}$  is such that  $p^{\alpha} \mid n$  but  $p^{\alpha+1} \nmid n$ . Deduce that there exists a positive constant  $C(\epsilon)$  depending only on  $\epsilon$  such that for all  $n \in \mathbb{N}$ ,  $\tau(n) \leq C(\epsilon)n^{\epsilon}$ .

# 12H Topics in Analysis

- (a) State Liouville's theorem on the approximation of algebraic numbers by rationals.
- (b) Let  $(a_n)_{n=0}^{\infty}$  be a sequence of positive integers and let

$$\alpha = a_0 + \frac{1}{a_1 + \frac{1}{a_2 + \frac{1}{a_3 + \dots}}}$$

be the value of the associated continued fraction.

(i) Prove that the *n*th convergent  $p_n/q_n$  satisfies

$$\left|\alpha - \frac{p_n}{q_n}\right| \leqslant \left|\alpha - \frac{p}{q}\right|$$

for all the rational numbers  $\frac{p}{q}$  such that  $0 < q \leq q_n$ .

(ii) Show that if the sequence  $(a_n)$  is bounded, then one can choose c > 0 (depending only on  $\alpha$ ), so that for every rational number  $\frac{a}{b}$ ,

$$\left|\alpha - \frac{a}{b}\right| > \frac{c}{b^2}$$

(iii) Show that if the sequence  $(a_n)$  is unbounded, then for each c > 0 there exist infinitely many rational numbers  $\frac{a}{b}$  such that

$$\left|\alpha - \frac{a}{b}\right| < \frac{c}{b^2}.$$

You may assume without proof the relation

$$\begin{pmatrix} p_{n+1} & p_n \\ q_{n+1} & q_n \end{pmatrix} = \begin{pmatrix} p_n & p_{n-1} \\ q_n & q_{n+1} \end{pmatrix} \begin{pmatrix} a_{n+1} & 1 \\ 1 & 0 \end{pmatrix}, \quad n = 1, 2, \dots$$

### 13J Statistical Modelling

(a) Define a generalised linear model (GLM) with design matrix  $X \in \mathbb{R}^{n \times p}$ , output variables  $Y := (Y_1, \ldots, Y_n)^T$  and parameters  $\mu := (\mu_1, \ldots, \mu_n)^T$ ,  $\beta \in \mathbb{R}^p$  and  $\sigma_i^2 := a_i \sigma^2 \in (0, \infty), i = 1, \ldots, n$ . Derive the moment generating function of Y, i.e. give an expression for  $\mathbb{E}\left[\exp\left(t^T Y\right)\right], t \in \mathbb{R}^n$ , wherever it is well-defined.

Assume from now on that the GLM satisfies the usual regularity assumptions, X has full column rank, and  $\sigma^2$  is known and satisfies  $1/\sigma^2 \in \mathbb{N}$ .

(b) Let  $\tilde{Y} := (\tilde{Y}_1, \dots, \tilde{Y}_{n/\sigma^2})^T$  be the output variables of a GLM from the same family as that of part (a) and parameters  $\tilde{\mu} := (\tilde{\mu}_1, \dots, \tilde{\mu}_{n/\sigma^2})^T$  and  $\tilde{\sigma}^2 := (\tilde{\sigma}_1^2, \dots, \tilde{\sigma}_{n/\sigma^2}^2)$ . Suppose the output variables may be split into n blocks of size  $1/\sigma^2$  with constant parameters. To be precise, for each block  $i = 1, \dots, n$ , if  $j \in \{(i-1)/\sigma^2 + 1, \dots, i/\sigma^2\}$  then

$$\tilde{\mu}_j = \mu_i$$
 and  $\tilde{\sigma}_i^2 = a_i$ 

with  $\mu_i = \mu_i(\beta)$  and  $a_i$  defined as in part (a). Let  $\bar{Y} := (\bar{Y}_1, \ldots, \bar{Y}_n)^T$ , where  $\bar{Y}_i := \sigma^2 \sum_{k=1}^{1/\sigma^2} \tilde{Y}_{(i-1)/\sigma^2+k}$ .

(i) Show that  $\overline{Y}$  is equal to Y in distribution. [*Hint: you may use without proof that moment generating functions uniquely determine distributions from exponential dispersion families.*]

(ii) For any  $\tilde{y} \in \mathbb{R}^{n/\sigma^2}$ , let  $\bar{y} = (\bar{y}_1, \dots, \bar{y}_n)^T$ , where  $\bar{y}_i := \sigma^2 \sum_{k=1}^{1/\sigma^2} \tilde{y}_{(i-1)/\sigma^2+k}$ . Show that the model function of  $\tilde{Y}$  satisfies

$$f\left(\tilde{y};\tilde{\mu},\tilde{\sigma}^{2}\right) = g_{1}\left(\bar{y};\tilde{\mu},\tilde{\sigma}^{2}\right) \times g_{2}\left(\tilde{y};\tilde{\sigma}^{2}\right)$$

for some functions  $g_1, g_2$ , and conclude that  $\overline{Y}$  is a sufficient statistic for  $\beta$  from  $\widetilde{Y}$ .

(iii) For the model and data from part (a), let  $\hat{\mu}$  be the maximum likelihood estimator for  $\mu$  and let  $D(Y;\mu)$  be the deviance at  $\mu$ . Using (i) and (ii), show that

$$\frac{D(Y;\hat{\mu})}{\sigma^2} =^d 2 \log \left\{ \frac{\sup_{\tilde{\mu}' \in \widetilde{\mathcal{M}}_1} f(\tilde{Y}; \tilde{\mu}', \tilde{\sigma}^2)}{\sup_{\tilde{\mu}' \in \widetilde{\mathcal{M}}_0} f(\tilde{Y}; \tilde{\mu}', \tilde{\sigma}^2)} \right\},$$

where  $=^d$  means equality in distribution and  $\widetilde{\mathcal{M}}_0$  and  $\widetilde{\mathcal{M}}_1$  are nested subspaces of  $\mathbb{R}^{n/\sigma^2}$  which you should specify. Argue that  $\dim(\widetilde{\mathcal{M}}_1) = n$  and  $\dim(\widetilde{\mathcal{M}}_0) = p$ , and, assuming the usual regularity assumptions, conclude that

$$\frac{D(Y;\hat{\mu})}{\sigma^2} \to^d \chi^2_{n-p} \qquad \text{as } \sigma^2 \to 0$$

stating the name of the result from class that you use.

### 14B Mathematical Biology

Consider the stochastic catalytic reaction

$$E \leftrightarrows ES, \qquad ES \rightarrow E + P$$

in which a single enzyme E complexes reversibly to ES (at forward rate  $k_1$  and reverse rate  $k'_1$ ) and decomposes into product P (at forward rate  $k_2$ ), regenerating enzyme E. Assume there is sufficient substrate S so that this catalytic cycle can continue indefinitely. Let P(E, n) be the probability of the state with enzyme E and n products and P(ES, n) the probability of the state with complex ES and n products, these states being mutually exclusive.

(i) Write down the master equation for the probabilities P(E, n) and P(ES, n) for  $n \ge 0$ .

(ii) Assuming an initial state with zero products, solve the master equation for P(E, 0) and P(ES, 0).

(iii) Hence find the probability distribution  $f(\tau)$  of the time  $\tau$  taken to form the first product.

(iv) Obtain the mean of  $\tau$ .

#### 15B Classical Dynamics

(a) Explain how the Hamiltonian  $H(\mathbf{q}, \mathbf{p}, t)$  of a system can be obtained from its Lagrangian  $L(\mathbf{q}, \dot{\mathbf{q}}, t)$ . Deduce that the action can be written as

$$S = \int (\mathbf{p} \cdot d\mathbf{q} - H \, dt) \, .$$

Show that Hamilton's equations are obtained if the action, computed between fixed initial and final configurations  $\mathbf{q}(t_1)$  and  $\mathbf{q}(t_2)$ , is minimized with respect to independent variations of  $\mathbf{q}$  and  $\mathbf{p}$ .

(b) Let  $(\mathbf{Q}, \mathbf{P})$  be a new set of coordinates on the same phase space. If the old and new coordinates are related by a type-2 generating function  $F_2(\mathbf{q}, \mathbf{P}, t)$  such that

$$\mathbf{p} = \frac{\partial F_2}{\partial \mathbf{q}} \,, \qquad \mathbf{Q} = \frac{\partial F_2}{\partial \mathbf{P}} \,,$$

deduce that the canonical form of Hamilton's equations applies in the new coordinates, but with a new Hamiltonian given by

$$K = H + \frac{\partial F_2}{\partial t} \,.$$

(c) For each of the Hamiltonians

(i) 
$$H = H(p)$$
, (ii)  $H = \frac{1}{2}(q^2 + p^2)$ ,

express the general solution (q(t), p(t)) at time t in terms of the initial values given by (Q, P) = (q(0), p(0)) at time t = 0. In each case, show that the transformation from (q, p) to (Q, P) is canonical for all values of t, and find the corresponding generating function  $F_2(q, P, t)$  explicitly.

Part II, Paper 4

## [TURN OVER]

# 16H Logic and Set Theory

(a) State Zorn's lemma.

[Throughout the remainder of this question, assume Zorn's lemma.]

(b) Let P be a poset in which every non-empty chain has an upper bound and let  $x \in P$ . By considering the poset  $P_x = \{y \in P \mid x \leq y\}$ , show that P has a maximal element  $\sigma$  with  $x \leq \sigma$ .

(c) A filter is a non-empty subset  $\mathcal{F} \subset \mathcal{P}(\mathbb{N})$  satisfying the following three conditions:

- if  $A, B \in \mathcal{F}$  then  $A \cap B \in \mathcal{F}$ ;
- if  $A \in \mathcal{F}$  and  $A \subset B$  then  $B \in \mathcal{F}$ ;
- $\emptyset \notin \mathcal{F}$ .

An ultrafilter is a filter  $\mathcal{U}$  such that for all  $A \subset \mathbb{N}$  we have either  $A \in \mathcal{U}$  or  $A^c \in \mathcal{U}$ , where  $A^c = \mathbb{N} \setminus A$ .

(i) For each  $n \in \mathbb{N}$ , show that  $\mathcal{U}_n = \{A \subset \mathbb{N} \mid n \in A\}$  is an ultrafilter.

(ii) Show that  $\mathcal{F} = \{A \subset \mathbb{N} \mid A^c \text{ is finite}\}$  is a filter but not an ultrafilter, and that for all  $n \in \mathbb{N}$  we have  $\mathcal{F} \not\subset \mathcal{U}_n$ .

(iii) Does there exist an ultrafilter  $\mathcal{U}$  such that  $\mathcal{U} \neq \mathcal{U}_n$  for any  $n \in \mathbb{N}$ ? Justify your answer.

#### 17G Graph Theory

State and prove Vizing's theorem about the chromatic index  $\chi'(G)$  of a graph G.

Let  $K_{m,n}$  be the complete bipartite graph with class sizes m and n. By first considering  $\chi'(K_{n,n})$ , find  $\chi'(K_{m,n})$  for all m and n.

Let G be the graph of order 2n + 1 obtained by subdividing a single edge of  $K_{n,n}$  by a new vertex. Show that  $\chi'(G) = \Delta(G) + 1$ , where  $\Delta(G)$  is the maximum degree of G.

## 18G Galois Theory

(a) Let K be a field. Define the discriminant  $\Delta(f)$  of a polynomial  $f(x) \in K[x]$ , and explain why it is in K, carefully stating any theorems you use.

Compute the discriminant of  $x^4 + rx + s$ .

(b) Let K be a field and let  $f(x) \in K[x]$  be a quartic polynomial with roots  $\alpha_1, \ldots, \alpha_4$  such that  $\alpha_1 + \cdots + \alpha_4 = 0$ .

Define the resolvant cubic g(x) of f(x).

Suppose that  $\Delta(f)$  is a square in K. Prove that the resolvant cubic is irreducible if and only if  $Gal(f) = A_4$ . Determine the possible Galois groups Gal(f) if g(x) is reducible.

The resolvant cubic of  $x^4 + rx + s$  is  $x^3 - 4sx - r^2$ . Using this, or otherwise, determine Gal(f), where  $f(x) = x^4 + 8x + 12 \in \mathbb{Q}[x]$ . [You may use without proof that f is irreducible.]

#### **19F** Representation Theory

(a) State and prove Burnside's lemma. Deduce that if a finite group G acts 2-transitively on a set X then the corresponding permutation character has precisely two (distinct) irreducible summands.

(b) Suppose that  $\mathbb{F}_q$  is a field with q elements. Write down a list of conjugacy class representatives for  $GL_2(\mathbb{F}_q)$ . Consider the natural action of  $GL_2(\mathbb{F}_q)$  on the set of lines through the origin in  $\mathbb{F}_q^2$ . What values does the corresponding permutation character take on each conjugacy class representative in your list? Decompose this permutation character into irreducible characters.

#### 20G Number Fields

Let K be a number field of degree n, and let  $\{\sigma_i \colon K \hookrightarrow \mathbb{C}\}$  be the set of complex embeddings of K. Show that if  $\alpha \in \mathcal{O}_K$  satisfies  $|\sigma_i(\alpha)| = 1$  for every i, then  $\alpha$  is a root of unity. Prove also that there exists c > 1 such that if  $\alpha \in \mathcal{O}_K$  and  $|\sigma_i(\alpha)| < c$  for all i, then  $\alpha$  is a root of unity.

State Dirichlet's Unit theorem.

Let  $K \subset \mathbb{R}$  be a real quadratic field. Assuming Dirichlet's Unit theorem, show that the set of units of K which are greater than 1 has a smallest element  $\epsilon$ , and that the group of units of Kis then  $\{\pm \epsilon^n \mid n \in \mathbb{Z}\}$ . Determine  $\epsilon$  for  $\mathbb{Q}(\sqrt{11})$ , justifying your result. [If you use the continued fraction algorithm, you must prove it in full.]

### 21F Algebraic Topology

In this question, you may assume all spaces involved are triangulable.

(a) (i) State and prove the Mayer–Vietoris theorem. [You may assume the theorem that states that a short exact sequence of chain complexes gives rise to a long exact sequence of homology groups.]

(ii) Use Mayer–Vietoris to calculate the homology groups of an oriented surface of genus g.

(b) Let S be an oriented surface of genus g, and let  $D_1, \ldots, D_n$  be a collection of mutually disjoint closed subsets of S with each  $D_i$  homeomorphic to a two-dimensional disk. Let  $D_i^{\circ}$  denote the interior of  $D_i$ , homeomorphic to an open two-dimensional disk, and let

$$T := S \setminus (D_1^{\circ} \cup \dots \cup D_n^{\circ}).$$

Show that

$$H_i(T) = \begin{cases} \mathbb{Z} & i = 0, \\ \mathbb{Z}^{2g+n-1} & i = 1, \\ 0 & \text{otherwise.} \end{cases}$$

(c) Let T be the surface given in (b) when  $S = S^2$  and n = 3. Let  $f: T \to S^1 \times S^1$  be a map. Does there exist a map  $g: S^1 \times S^1 \to T$  such that  $f \circ g$  is homotopic to the identity map? Justify your answer.

#### 22I Linear Analysis

(a) For K a compact Hausdorff space, what does it mean to say that a set  $S \subset C(K)$  is equicontinuous. State and prove the Arzelà-Ascoli theorem.

(b) Suppose K is a compact Hausdorff space for which C(K) is a countable union of equicontinuous sets. Prove that K is finite.

(c) Let  $F : \mathbb{R}^n \to \mathbb{R}^n$  be a bounded, continuous function and let  $x_0 \in \mathbb{R}^n$ . Consider the problem of finding a differentiable function  $x : [0, 1] \to \mathbb{R}^n$  with

$$x(0) = x_0$$
 and  $x'(t) = F(x(t))$  (\*)

for all  $t \in [0,1]$ . For each  $k = 1, 2, 3, \ldots$ , let  $x_k : [0,1] \to \mathbb{R}^n$  be defined by setting  $x_k(0) = x_0$  and

$$x_k(t) = x_0 + \int_0^t F(y_k(s)) \, ds$$

for  $t \in [0, 1]$ , where

$$y_k(t) = x_k\left(\frac{j}{k}\right)$$

for  $t \in (\frac{j}{k}, \frac{j+1}{k}]$  and  $j \in \{0, 1, \dots, k-1\}.$ 

(i) Verify that  $x_k$  is well-defined and continuous on [0, 1] for each k.

(ii) Prove that there exists a differentiable function  $x: [0,1] \to \mathbb{R}^n$  solving (\*) for  $t \in [0,1]$ .

# 23I Analysis of Functions

(a) Define the Sobolev space  $H^s(\mathbb{R}^n)$  for  $s \in \mathbb{R}$ .

(b) Let k be a non-negative integer and let  $s > k + \frac{n}{2}$ . Show that if  $u \in H^s(\mathbb{R}^n)$  then there exists  $u^* \in C^k(\mathbb{R}^n)$  with  $u = u^*$  almost everywhere.

(c) Show that if  $f \in H^s(\mathbb{R}^n)$  for some  $s \in \mathbb{R}$ , there exists a unique  $u \in H^{s+4}(\mathbb{R}^n)$  which solves:

$$\Delta \Delta u + \Delta u + u = f,$$

in a distributional sense. Prove that there exists a constant C > 0, independent of f, such that:

$$||u||_{H^{s+4}} \leq C ||f||_{H^s}.$$

For which s will u be a classical solution?

24F Algebraic Geometry

Let  $P_0, \ldots, P_n$  be a basis for the homogeneous polynomials of degree n in variables  $Z_0$  and  $Z_1$ . Then the image of the map  $\mathbb{P}^1 \to \mathbb{P}^n$  given by

$$[Z_0, Z_1] \mapsto [P_0(Z_0, Z_1), \dots, P_n(Z_0, Z_1)]$$

is called a rational normal curve.

Let  $p_1, \ldots, p_{n+3}$  be a collection of points in general linear position in  $\mathbb{P}^n$ . Prove that there exists a unique rational normal curve in  $\mathbb{P}^n$  passing through these points.

Choose a basis of homogeneous polynomials of degree 3 as above, and give generators for the homogeneous ideal of the corresponding rational normal curve.

#### 25I Differential Geometry

(a) State the Gauss–Bonnet theorem for compact regular surfaces  $S \subset \mathbb{R}^3$  without boundary. Identify all expressions occurring in any formulae.

(b) Let  $S \subset \mathbb{R}^3$  be a compact regular surface without boundary and suppose that its Gaussian curvature  $K(x) \ge 0$  for all  $x \in S$ . Show that S is diffeomorphic to the sphere.

Let  $S_n$  be a sequence of compact regular surfaces in  $\mathbb{R}^3$  and let  $K_n(x)$  denote the Gaussian curvature of  $S_n$  at  $x \in S_n$ . Suppose that

$$\limsup_{n \to \infty} \inf_{x \in S_n} K_n(x) \ge 0. \tag{(\star)}$$

(c) Give an example to show that it does not follow that for all sufficiently large n the surface  $S_n$  is diffeomorphic to the sphere.

(d) Now assume, in addition to  $(\star)$ , that all of the following conditions hold:

- (1) There exists a constant  $R < \infty$  such that for all  $n, S_n$  is contained in a ball of radius R around the origin.
- (2) There exists a constant  $M < \infty$  such that  $\operatorname{Area}(S_n) \leq M$  for all n.
- (3) There exists a constant  $\epsilon_0 > 0$  such that for all n, all points  $p \in S_n$  admit a geodesic polar coordinate system centred at p of radius at least  $\epsilon_0$ .
- (4) There exists a constant  $C < \infty$  such that on all such geodesic polar neighbourhoods,  $|\partial_r K_n| \leq C$  for all n, where r denotes a geodesic polar coordinate.

(i) Show that for all sufficiently large n, the surface  $S_n$  is diffeomorphic to the sphere. [Hint: It may be useful to identify a geodesic polar ball  $B(p_n, \epsilon_0)$  in each  $S_n$  for which  $\int_{B(p_n, \epsilon_0)} K_n dA$  is bounded below by a positive constant independent of n.]

(ii) Explain how your example from (c) fails to satisfy one or more of these extra conditions (1)-(4).

[You may use without proof the standard computations for geodesic polar coordinates: E = 1, F = 0,  $\lim_{r\to 0} G(r, \theta) = 0$ ,  $\lim_{r\to 0} (\sqrt{G})_r(r, \theta) = 1$ , and  $(\sqrt{G})_{rr} = -K\sqrt{G}$ .]

Part II, Paper 4

#### 26K Probability and Measure

(a) State and prove the strong law of large numbers for sequences of i.i.d. random variables with a finite moment of order 4.

(b) Let  $(X_k)_{k \ge 1}$  be a sequence of independent random variables such that

$$\mathbb{P}(X_k = 1) = \mathbb{P}(X_k = -1) = \frac{1}{2}.$$

Let  $(a_k)_{k \ge 1}$  be a sequence of real numbers such that

$$\sum_{k \geqslant 1} a_k^2 < \infty.$$

Set

$$S_n := \sum_{k=1}^n a_k X_k.$$

(i) Show that  $S_n$  converges in  $\mathbb{L}^2$  to a random variable S as  $n \to \infty$ . Does it converge in  $\mathbb{L}^1$ ? Does it converge in law?

(ii) Show that  $||S||_4 \leq 3^{1/4} ||S||_2$ .

(iii) Let  $(Y_k)_{k \ge 1}$  be a sequence of i.i.d. standard Gaussian random variables, i.e. each  $Y_k$  is distributed as  $\mathcal{N}(0, 1)$ . Show that then  $\sum_{k=1}^n a_k Y_k$  converges in law as  $n \to \infty$  to a random variable and determine the law of the limit.

#### 27K Applied Probability

(i) Explain the notation  $M(\lambda)/M(\mu)/1$  in the context of queueing theory. [In the following, you may use without proof the fact that  $\pi_n = (\lambda/\mu)^n$  is the invariant distribution of such a queue when  $0 < \lambda < \mu$ .]

(ii) In a shop queue, some customers rejoin the queue after having been served. Let  $\lambda, \beta \in (0, \infty)$  and  $\delta \in (0, 1)$ . Consider a  $M(\lambda)/M(\mu)/1$  queue subject to the modification that, on completion of service, each customer leaves the shop with probability  $\delta$ , or rejoins the shop queue with probability  $1 - \delta$ . Different customers behave independently of one another, and all service times are independent random variables.

Find the distribution of the total time a given customer spends being served by the server. Hence show that equilibrium is possible if  $\lambda < \delta \mu$ , and find the invariant distribution of the queue-length in this case.

(iii) Show that, in equilibrium, the departure process is Poissonian, whereas, assuming the rejoining customers go to the end of the queue, the process of customers arriving at the queue (including the rejoining ones) is not Poissonian.

#### 28J Principles of Statistics

Consider  $X_1, \ldots, X_n$  drawn from a statistical model  $\{f(\cdot, \theta) : \theta \in \Theta\}, \Theta = \mathbb{R}^p$ , with nonsingular Fisher information matrix  $I(\theta)$ . For  $\theta_0 \in \Theta, h \in \mathbb{R}^p$ , define likelihood ratios

$$Z_n(h) = \log \frac{\prod_{i=1}^n f(X_i, \theta_0 + h/\sqrt{n})}{\prod_{i=1}^n f(X_i, \theta_0)}, \quad X_i \sim^{i.i.d.} f(\cdot, \theta_0).$$

Next consider the probability density functions  $(p_h : h \in \mathbb{R}^p)$  of normal distributions  $N(h, I(\theta_0)^{-1})$  with corresponding likelihood ratios given by

$$Z(h) = \log \frac{p_h(X)}{p_0(X)}, \ X \sim p_0.$$

Show that for every fixed  $h \in \mathbb{R}^p$ , the random variables  $Z_n(h)$  converge in distribution as  $n \to \infty$  to Z(h).

[You may assume suitable regularity conditions of the model  $\{f(\cdot, \theta) : \theta \in \Theta\}$  without specification, and results on uniform laws of large numbers from lectures can be used without proof.]

#### 29K Stochastic Financial Models

(i) What does it mean to say that  $(S_t^0, S_t)_{0 \leq t \leq T}$  is a *Black–Scholes model* with interest rate r, drift  $\mu$  and volatility  $\sigma$ ?

(ii) Write down the Black–Scholes pricing formula for the time-0 value  $V_0$  of a time-T contingent claim C.

(iii) Show that if C is a European call of strike K and maturity T then

$$V_0 \geqslant S_0 - e^{-rT} K.$$

(iv) For the European call, derive the Black–Scholes pricing formula

$$V_0 = S_0 \Phi(d^+) - e^{-rT} K \Phi(d^-),$$

where  $\Phi$  is the standard normal distribution function and  $d^+$  and  $d^-$  are to be determined.

(v) Fix  $t \in (0, T)$  and consider a modified contract which gives the investor the right but not the obligation to buy one unit of the risky asset at price K, either at time t or time T but not both, where the choice of exercise time is to be made by the investor at time t. Determine whether the investor should exercise the contract at time t.

#### **30J** Mathematics of Machine Learning

Suppose we have input-output pairs  $(x_1, y_1), \ldots, (x_n, y_n) \in \mathbb{R}^p \times \{-1, 1\}$ . Consider the empirical risk minimisation problem with hypothesis class

$$\mathcal{H} = \{ x \mapsto x^T \beta : \beta \in C \},\$$

where C is a non-empty closed convex subset of  $\mathbb{R}^p$ , and logistic loss

$$\ell(h(x), y) = \log_2(1 + e^{-yh(x)}),$$

for  $h \in \mathcal{H}$  and  $(x, y) \in \mathbb{R}^p \times \{-1, 1\}$ .

(i) Show that the objective function f of the optimisation problem is convex.

(ii) Let  $\pi_C(x)$  denote the projection of x onto C. Describe the procedure of *stochastic* gradient descent (SGD) for minimisation of f above, giving explicit forms for any gradients used in the algorithm.

(iii) Suppose  $\hat{\beta}$  minimises  $f(\beta)$  over  $\beta \in C$ . Suppose  $\max_{i=1,...,n} ||x_i||_2 \leq M$  and  $\sup_{\beta \in C} ||\beta||_2 \leq R$ . Prove that the output  $\bar{\beta}$  of k iterations of the SGD algorithm with some fixed step size  $\eta$  (which you should specify), satisfies

$$\mathbb{E}f(\bar{\beta}) - f(\hat{\beta}) \leqslant \frac{2MR}{\log(2)\sqrt{k}}.$$

(iv) Now suppose that the step size at iteration s is  $\eta_s > 0$  for each s = 1, ..., k. Show that, writing  $\beta_s$  for the sth iterate of SGD, we have

$$\mathbb{E}f(\tilde{\beta}) - f(\hat{\beta}) \leqslant \frac{A_2 M^2}{2A_1 \{\log(2)\}^2} + \frac{2R^2}{A_1},$$

where

$$\tilde{\beta} = \frac{1}{A_1} \sum_{s=1}^k \eta_s \beta_s, \quad A_1 = \sum_{s=1}^k \eta_s \quad \text{and} \quad A_2 = \sum_{s=1}^k \eta_s^2.$$

[You may use standard properties of convex functions and projections onto closed convex sets without proof provided you state them.]

# 31A Asymptotic Methods

Consider the differential equation

$$y'' - y' - \frac{2(x+1)}{x^2} y = 0.$$
<sup>(†)</sup>

(i) Classify what type of regularity/singularity equation (†) has at  $x = \infty$ .

(ii) Find a transformation that maps equation (†) to an equation of the form

$$u'' + q(x)u = 0. (*)$$

(iii) Find the leading-order term of the asymptotic expansions of the solutions of equation (\*), as  $x \to \infty$ , using the Liouville–Green method.

(iv) Derive the leading-order term of the asymptotic expansion of the solutions y of  $(\dagger)$ . Check that one of them is an exact solution for  $(\dagger)$ .

#### 32E Dynamical Systems

(a) Let  $F: I \to I$  be a continuous map defined on an interval  $I \subset \mathbb{R}$ . Define what it means (i) for F to have a *horseshoe* and (ii) for F to be *chaotic*. [Glendinning's definition should be used throughout this question.]

(b) Consider the map defined on the interval [-1, 1] by

$$F(x) = 1 - \mu |x|$$

with  $0 < \mu \leq 2$ .

- (i) Sketch F(x) and  $F^2(x)$  for a case when  $0 < \mu < 1$  and a case when  $1 < \mu < 2$ .
- (ii) Describe fully the long term dynamics for  $0 < \mu < 1$ . What happens for  $\mu = 1$ ?
- (iii) When does F have a horseshoe? When does  $F^2$  have a horseshoe?
- (iv) For what values of  $\mu$  is the map F chaotic?

#### 33A Principles of Quantum Mechanics

Briefly explain why the density operator  $\rho$  obeys  $\rho \ge 0$  and  $\text{Tr}(\rho) = 1$ . What is meant by a *pure* state and a *mixed* state?

A two-state system evolves under the Hamiltonian  $H = \hbar \omega \cdot \sigma$ , where  $\omega$  is a constant vector and  $\sigma$  are the Pauli matrices. At time t the system is described by a density operator

$$\rho(t) = \frac{1}{2} \left( 1_{\mathcal{H}} + \mathbf{a}(t) \cdot \boldsymbol{\sigma} \right)$$

where  $1_{\mathcal{H}}$  is the identity operator. Initially, the vector  $\mathbf{a}(0) = \mathbf{a}$  obeys  $|\mathbf{a}| < 1$  and  $\mathbf{a} \cdot \boldsymbol{\omega} = 0$ . Find  $\rho(t)$  in terms of  $\mathbf{a}$  and  $\boldsymbol{\omega}$ . At what time, if any, is the system definitely in the state  $|\uparrow_x\rangle$  that obeys  $\sigma_x|\uparrow_x\rangle = +|\uparrow_x\rangle$ ?

Part II, Paper 4

# [TURN OVER]

# 34C Applications of Quantum Mechanics

(a) For a particle of charge q moving in an electromagnetic field with vector potential  $\mathbf{A}$  and scalar potential  $\phi$ , write down the classical Hamiltonian and the equations of motion.

(b) Consider the vector and scalar potentials

$$\boldsymbol{A} = \frac{B}{2}(-y, x, 0), \qquad \phi = 0$$

(i) Solve the equations of motion. Define and compute the cyclotron frequency  $\omega_B$ .

(ii) Write down the quantum Hamiltonian of the system in terms of the angular momentum operator

$$L_z = xp_y - yp_x$$

Show that the states

$$\psi(x,y) = f(x+iy)e^{-(x^2+y^2)qB/4\hbar},$$
(†)

for any function f, are energy eigenstates and compute their energy. Define Landau levels and discuss this result in relation to them.

(iii) Show that for  $f(w) = w^M$ , the wavefunctions in (†) are eigenstates of angular momentum and compute the corresponding eigenvalue. These wavefunctions peak in a ring around the origin. Estimate its radius. Using these two facts or otherwise, estimate the degeneracy of Landau levels.

# 35A Statistical Physics

Consider a classical gas of N particles in volume V, where the total energy is the standard kinetic energy plus a potential  $U(\mathbf{x}_1, \mathbf{x}_2, ..., \mathbf{x}_N)$  depending on the relative locations of the particles  $\{\mathbf{x}_i : 1 \leq i \leq N\}$ .

(i) Starting from the partition function, show that the free energy of the gas is

$$F = F_{\text{ideal}} - T \log \left\{ 1 + \frac{1}{V^N} \int (e^{-U/T} - 1) d^{3N} x \right\}, \qquad (*)$$

where  $F_{\text{ideal}}$  is the free energy when  $U \equiv 0$ .

(ii) Suppose now that the gas is fairly dilute and that the integral in (\*) is small compared to  $V^N$  and is dominated by two-particle interactions. Show that the free energy simplifies to the form

$$F = F_{\text{ideal}} + \frac{N^2 T}{V} B(T), \tag{\dagger}$$

and find an integral expression for B(T). Using (†) find the equation of state of the gas, and verify that B(T) is the second virial coefficient.

(iii) The equation of state for a Clausius gas is

$$P(V - Nb) = NT$$

for some constant b. Find the second virial coefficient for this gas. Evaluate b for a gas of hard sphere atoms of radius  $r_0$ .

#### 36D Electrodynamics

(a) A dielectric medium exhibits a linear response if the electric displacement  $\mathbf{D}(\mathbf{x}, t)$  and magnetizing field  $\mathbf{H}(\mathbf{x}, t)$  are related to the electric and magnetic fields,  $\mathbf{E}(\mathbf{x}, t)$  and  $\mathbf{B}(\mathbf{x}, t)$ , as

$$\mathbf{D} = \epsilon \mathbf{E} \,, \qquad \mathbf{B} = \mu \mathbf{H} \,,$$

where  $\epsilon$  and  $\mu$  are constants characterising the electric and magnetic polarisability of the material respectively. Write down the Maxwell equations obeyed by the fields **D**, **H**, **B** and **E** in this medium in the absence of free charges or currents.

(b) Two such media with constants  $\epsilon_{-}$  and  $\epsilon_{+}$  (but the same  $\mu$ ) fill the regions x < 0 and x > 0 respectively in three-dimensions with Cartesian coordinates (x, y, z).

(i) Starting from Maxwell's equations, derive the appropriate boundary conditions at x = 0 for a time-independent electric field  $\mathbf{E}(\mathbf{x})$ .

(ii) Consider a candidate solution of Maxwell's equations describing the reflection and transmission of an incident electromagnetic wave of wave vector  $\mathbf{k}_{I}$  and angular frequency  $\omega_{I}$  off the interface at x = 0. The electric field is given as,

$$\mathbf{E}(\mathbf{x},t) = \begin{cases} \sum_{X=I,R} \operatorname{Im}\left[\mathbf{E}_{X} \exp\left(i\mathbf{k}_{X}\cdot\mathbf{x} - i\omega_{X}t\right)\right], & x < 0, \\ \operatorname{Im}\left[\mathbf{E}_{T} \exp\left(i\mathbf{k}_{T}\cdot\mathbf{x} - i\omega_{T}t\right)\right], & x > 0, \end{cases}$$

where  $\mathbf{E}_I$ ,  $\mathbf{E}_R$  and  $\mathbf{E}_T$  are constant real vectors and Im[z] denotes the imaginary part of a complex number z. Give conditions on the parameters  $\mathbf{E}_X$ ,  $\mathbf{k}_X$ ,  $\omega_X$  for X = I, R, T, such that the above expression for the electric field  $\mathbf{E}(\mathbf{x}, t)$  solves Maxwell's equations for all  $x \neq 0$ , together with an appropriate magnetic field  $\mathbf{B}(\mathbf{x}, t)$  which you should determine.

(iii) We now parametrize the incident wave vector as  $\mathbf{k}_I = k_I(\cos(\theta_I)\hat{\mathbf{i}}_x + \sin(\theta_I)\hat{\mathbf{i}}_z)$ , where  $\hat{\mathbf{i}}_x$  and  $\hat{\mathbf{i}}_z$  are unit vectors in the x- and z-directions respectively, and choose the incident polarisation vector to satisfy  $\mathbf{E}_I \cdot \hat{\mathbf{i}}_x = 0$ . By imposing appropriate boundary conditions for  $\mathbf{E}(\mathbf{x}, t)$  at x = 0, which you may assume to be the same as those for the time-independent case considered above, determine the Cartesian components of the wavevector  $\mathbf{k}_T$  as functions of  $k_I$ ,  $\theta_I$ ,  $\epsilon_+$  and  $\epsilon_-$ .

(iv) For  $\epsilon_+ < \epsilon_-$  find a critical value  $\theta_I^{cr}$  of the angle of incidence  $\theta_I$  above which there is no real solution for the wavevector  $\mathbf{k}_T$ . Write down a solution for  $\mathbf{E}(\mathbf{x}, t)$  when  $\theta_I > \theta_I^{cr}$  and comment on its form.

#### 37D General Relativity

In linearized general relativity, we consider spacetime metrics that are perturbatively close to Minkowski,  $g_{\mu\nu} = \eta_{\mu\nu} + h_{\mu\nu}$ , where  $\eta_{\mu\nu} = \text{diag}(-1, 1, 1, 1)$  and  $h_{\mu\nu} = \mathcal{O}(\epsilon) \ll 1$ . In the Lorenz gauge, the Einstein tensor, at linear order, is given by

$$G_{\mu\nu} = -\frac{1}{2} \Box \bar{h}_{\mu\nu} , \qquad \bar{h}_{\mu\nu} = h_{\mu\nu} - \frac{1}{2} \eta_{\mu\nu} h , \qquad (\dagger)$$

where  $\Box = \eta^{\mu\nu} \partial_{\mu} \partial_{\nu}$  and  $h = \eta^{\mu\nu} h_{\mu\nu}$ .

(i) Show that the (fully nonlinear) Einstein equations  $G_{\alpha\beta} = 8\pi T_{\alpha\beta}$  can be equivalently written in terms of the Ricci tensor  $R_{\alpha\beta}$  as

$$R_{\alpha\beta} = 8\pi \left( T_{\alpha\beta} - \frac{1}{2} g_{\alpha\beta} T \right) , \qquad T = g^{\mu\nu} T_{\mu\nu} .$$

Show likewise that equation  $(\dagger)$  can be written as

$$\Box h_{\mu\nu} = -16\pi \left( T_{\mu\nu} - \frac{1}{2} \eta_{\mu\nu} T \right) \,. \tag{(*)}$$

(ii) In the Newtonian limit we consider matter sources with small velocities  $v \ll 1$  such that time derivatives  $\partial/\partial t \sim v \partial/\partial x^i$  can be neglected relative to spatial derivatives, and the only non-negligible component of the energy-momentum tensor is the energy density  $T_{00} = \rho$ . Show that in this limit, we recover from equation (\*) the Poisson equation  $\vec{\nabla}^2 \Phi = 4\pi\rho$  of Newtonian gravity if we identify  $h_{00} = -2\Phi$ .

(iii) A point particle of mass M is modelled by the energy density  $\rho = M \,\delta(\mathbf{r})$ . Derive the Newtonian potential  $\Phi$  for this point particle by solving the Poisson equation.

[You can assume the solution of 
$$\vec{\nabla}^2 \varphi = f(\mathbf{r})$$
 is  $\varphi(\mathbf{r}) = -\int \frac{f(\mathbf{r}')}{4\pi |\mathbf{r} - \mathbf{r}'|} d^3 \mathbf{r}'$ .]

(iv) Now consider the Einstein equations with a small positive cosmological constant,  $G_{\alpha\beta} + \Lambda g_{\alpha\beta} = 8\pi T_{\alpha\beta}, \Lambda = \mathcal{O}(\epsilon) > 0$ . Repeat the steps of questions (i)-(iii), again identifying  $h_{00} = -2\Phi$ , to show that the Newtonian limit is now described by the Poisson equation  $\vec{\nabla}^2 \Phi = 4\pi\rho - \Lambda$ , and that a solution for the potential of a point particle is given by

$$\Phi = -\frac{M}{r} - Br^2 \,,$$

where *B* is a constant you should determine. Briefly discuss the effect of the  $Br^2$  term and determine for which range of the radius *r* the weak-field limit is a justified approximation. [*Hint:* Absorb the term  $\Lambda g_{\alpha\beta}$  as part of the energy-momentum tensor. Note also that in spherical symmetry  $\vec{\nabla}^2 f = \frac{1}{r} \frac{\partial^2}{\partial r^2} (rf)$ .]

# 38B Fluid Dynamics II

Consider a two-dimensional horizontal vortex sheet of strength U in a homogeneous inviscid fluid at height h above a horizontal rigid boundary at y = 0 so that the fluid velocity is

$$\boldsymbol{u} = \left\{ \begin{array}{ll} U \hat{\boldsymbol{x}} \,, & 0 < y < h \,, \\ \boldsymbol{0} \,, & h < y \,. \end{array} \right.$$

(i) Investigate the linear instability of the sheet by determining the relevant dispersion relation for small, inviscid, irrotational perturbations. For what wavelengths is the sheet unstable?

(ii) Evaluate the temporal growth rate and the wave propagation speed in the limits of both short and long waves. Using these results, sketch how the growth rate varies with the wavenumber.

(iii) Comment briefly on how the introduction of a stable density difference (fluid in y > h is less dense than that in 0 < y < h) and surface tension at the interface would affect the growth rates.

#### 39B Waves

(a) Show that the equations for one-dimensional unsteady flow of an inviscid compressible fluid at constant entropy can be put in the form

$$\Big(\frac{\partial}{\partial t} + (u \pm c)\frac{\partial}{\partial x}\Big)R_{\pm} = 0\,,$$

where u and c are the fluid velocity and the local sound speed, respectively, and the Riemann invariants  $R_{\pm}$  are to be defined.

Such a fluid occupies a long narrow tube along the x-axis. For times t < 0 it is at rest with uniform pressure  $p_0$ , density  $\rho_0$  and sound speed  $c_0$ . At t = 0 a finite segment,  $0 \le x \le L$ , is disturbed so that u = U(x) and  $c = c_0 + C(x)$ , with U = C = 0 for  $x \le 0$  and  $x \ge L$ . Explain, with the aid of a carefully labelled sketch, how two independent simple waves emerge after some time. You may assume that no shock waves form.

(b) A fluid has the adiabatic equation of state

$$p(\rho) = A - \frac{B^2}{\rho} \,,$$

where A and B are positive constants and  $\rho > B^2/A$ .

(i) Calculate the Riemann invariants for this fluid, and express  $u \pm c$  in terms of  $R_{\pm}$  and  $c_0$ . Deduce that in a simple wave with  $R_{-} = 0$  the velocity field translates, without any nonlinear distortion, at the equilibrium sound speed  $c_0$ .

(ii) At t = 0 this fluid occupies x > 0 and is at rest with uniform pressure, density and sound speed. For t > 0 a piston initially at x = 0 executes simple harmonic motion with position  $x(t) = a \sin \omega t$ , where  $a\omega < c_0$ . Show that  $u(x,t) = U(\phi)$ , where  $\phi = \omega(t-x/c_0)$ , for some function U that is zero for  $\phi < 0$  and is  $2\pi$ -periodic, but not simple harmonic, for  $\phi > 0$ . By approximately inverting the relationship between  $\phi$  and the time  $\tau$  that a characteristic leaves the piston for the case  $\epsilon = a\omega/c_0 \ll 1$ , show that

$$U(\phi) = a\omega \left(\cos\phi - \epsilon \sin^2\phi - \frac{3}{2}\epsilon^2 \sin^2\phi \cos\phi + O(\epsilon^3)\right) \quad \text{for} \quad \phi > 0 \,.$$

## 40E Numerical Analysis

(a) For a function f = f(x, y) which is real analytic in  $\mathbb{R}^2$  and 2-periodic in each variable, its Fourier expansion is given by the formula

$$f(x,y) = \sum_{m,n\in\mathbb{Z}} \widehat{f}_{m,n} e^{i\pi mx + i\pi ny}, \qquad \widehat{f}_{m,n} = \frac{1}{4} \int_{-1}^{1} \int_{-1}^{1} f(t,\theta) e^{-i\pi mt - i\pi n\theta} dt d\theta.$$

Derive expressions for the Fourier coefficients of partial derivatives  $f_x$ ,  $f_y$  and those of the product h(x, y) = f(x, y)g(x, y) in terms of  $\hat{f}_{m,n}$  and  $\hat{g}_{m,n}$ .

(b) Let u(x,y) be the 2-periodic solution in  $\mathbb{R}^2$  of the general second-order elliptic PDE

$$(au_x)_x + (au_y)_y = f,$$

where a and f are both real analytic and 2-periodic, and a(x, y) > 0. We impose the normalisation condition  $\int_{-1}^{1} \int_{-1}^{1} u \, dx \, dy = 0$  and note from the PDE  $\int_{-1}^{1} \int_{-1}^{1} f \, dx \, dy = 0$ .

Construct explicitly the infinite-dimensional linear algebraic system that arises from the application of the Fourier spectral method to the above equation, and explain how to truncate this system to a finite-dimensional one.

(c) Specify the truncated system for the unknowns  $\{\hat{u}_{m,n}\}\$  for the case

$$a(x,y) = 5 + 2\cos\pi x + 2\cos\pi y,$$

and prove that, for any ordering of the Fourier coefficients  $\{\hat{u}_{m,n}\}$  into one-dimensional array, the resulting system is symmetric and positive definite. [You may use the Gershgorin theorem without proof.]

# END OF PAPER