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## United Kingdom

Short-listed Problems and Solutions

## Short-listed Problems and Solutions

The Problem Selection Committee received a total of 130 problem proposals, from 41 countries. In this booklet, a short-list of 27 problems is presented, classified under Number Theory (N1N6), Geometry (G1-G8), Algebra (A1-A6) and Combinatorics (C1-C7). Within each section, the problems are arranged in ascending order of estimated difficulty, although of course it is very hard to judge this accurately.

The Problem Selection Committee would like to thank the following countries for submitting problems:

| Albania | Greece | Poland |
| :--- | :--- | :--- |
| Armenia | Holland | Puerto Rico |
| Australia | Hong Kong | Romania |
| Brazil | India | Russia |
| Bulgaria | Iran | Slovakia |
| Canada | Ireland | Slovenia |
| Colombia | Israel | Sweden |
| Croatia | Japan | Taiwan (ROC) |
| Czech Republic | Korea | Thailand |
| Estonia | Lithuania | Ukraine |
| Finland | Luxembourg | United States |
| France | Moldova | Uzbekistan |
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| N3 3 | G3 | 11 | A3 | 24 | C3 | 35 |
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Typeset by Bill Richardson.

N1. What is the smallest positive integer $t$ such that there exist integers $x_{1}, x_{2}, \ldots, x_{t}$ with

$$
x_{1}^{3}+x_{2}^{3}+\ldots+x_{t}^{3}=2002^{2002} ?
$$

Solution. The answer is $t=4$.
We first show that $2002^{2002}$ is not a sum of three cubes by considering numbers modulo 9 . Thus, from $2002 \equiv 4(\bmod 9), 4^{3} \equiv 1(\bmod 9)$ and $2002=667 \times 3+1$ we find that

$$
2002^{2002} \equiv 4^{2002} \equiv 4(\bmod 9)
$$

whereas, from $x^{3} \equiv 0, \pm 1(\bmod 9)$ for any integer $x$, we see that $x_{1}^{3}+x_{2}^{3}+x_{3}^{3} \not \equiv 4(\bmod 9)$. It remains to show that $2002^{2002}$ is a sum of four cubes. Starting with

$$
2002=10^{3}+10^{3}+1^{3}+1^{3}
$$

and using $2002=667 \times 3+1$ once again, we find that

$$
\begin{aligned}
2002^{2002} & =2002 \times\left(2002^{667}\right)^{3} \\
& =\left(10 \times 2002^{667}\right)^{3}+\left(10 \times 2002^{667}\right)^{3}+\left(2002^{667}\right)^{3}+\left(2002^{667}\right)^{3} .
\end{aligned}
$$

## Comments

1. This is an easy question. The only subtle point is that, to show that $2002^{2002}$ is not the sum of three cubes, we need to consider a non-prime modulus. Indeed, to restrict the number of cubes $\bmod n$ we would like $\phi(n)$ to be a multiple of 3 (so that Fermat-Euler is helping us), but taking $n$ to be 7 or 13 or 19 does not help: there are too many cubes. So we try a composite $n$ with $\phi(n)$ a multiple of 3 , and the first such is $n=9$.
2. The proposer's original version of the problem only asked for a proof that three cubes is impossible and five cubes is possible. It is a fortunate feature of the number $2002^{2002}$ that we are able to settle the case of four cubes.

N2. Let $n \geqslant 2$ be a positive integer, with divisors $1=d_{1}<d_{2}<\ldots<d_{k}=n$. Prove that $d_{1} d_{2}+d_{2} d_{3}+\ldots+d_{k-1} d_{k}$ is always less than $n^{2}$, and determine when it is a divisor of $n^{2}$.

Solution. Note that if $d$ is a divisor of $n$ then so is $n / d$, so that the sum

$$
s=\sum_{1 \leqslant i<k} d_{i} d_{i+1}=n^{2} \sum_{1 \leqslant i<k} \frac{1}{d_{i} d_{i+1}} \leqslant n^{2} \sum_{1 \leqslant i<k}\left(\frac{1}{d_{i}}-\frac{1}{d_{i+1}}\right)<\frac{n^{2}}{d_{1}}=n^{2}
$$

Note also that $d_{2}=p, d_{k-1}=n / p, d_{k}=n$, where $p$ is the least prime divisor of $n$.
If $n=p$ then $k=2$ and $s=p$, which divides $n^{2}$.
If $n$ is composite then $k>2$, and $s>d_{k-1} d_{k}=n^{2} / p$. If such an $s$ were a divisor of $n^{2}$ then also $n^{2} / s$ would be a divisor of $n^{2}$. But $1<n^{2} / s<p$, which is impossible because $p$ is the least prime divisor of $n^{2}$.
Hence, the given sum is a divisor of $n^{2}$ if and only if $n$ is prime.

## Comments

1. The problem is perhaps not quite as easy as the short solution here appears to suggest. Even having done the first part, it is very easy to get stuck on the second part.
2. It would be possible to delete from the question the fact that the given expression is always less than $n^{2}$. But, in our opinion, the form as given above is natural and inviting to a reader.

N3. Let $p_{1}, p_{2}, \ldots, p_{n}$ be distinct primes greater than 3 . Show that $2^{p_{1} p_{2} \ldots p_{n}}+1$ has at least $4^{n}$ divisors.

## Comment

1. The natural strategy for this problem is to use induction on the number of primes involved, hoping that the number of divisors increases by a factor of 4 for each new prime in the expression. By the usual properties of the divisor function $d(m)$, it would be enough to show that $2^{p_{1} p_{2} \ldots p_{n}}+1$ contains at least two new prime factors not contained in $2^{p_{1} p_{2} \ldots p_{n-1}}+1$. Unfortunately this does not seem to be easy. Instead, we will show in an elementary way that there is at least one new prime at each step. To finish the proof, we will need the following additional observation: if $k>m$ then $d(k m) \geqslant 2 d(m)$, which follows from the simple fact that if $a$ divides $m$ then both $a$ and $k a$ divide $k m$.

Solution. We claim first that if $u$ and $v$ are coprime odd numbers then the highest common factor of $2^{u}+1$ and $2^{v}+1$ is 3 . Certainly 3 divides $2^{u}+1$ and $2^{v}+1$, because $u$ and $v$ are odd. Suppose now that some $t>3$ divides $2^{u}+1$ and $2^{v}+1$. Then we have $2^{u} \equiv-1(\bmod t)$ and $2^{v} \equiv-1(\bmod t)$. But if any $2^{x}$ is $-1 \bmod t$ then the set of all such $x$ is the set of all odd multiples of $r / 2$, where $r$ is the order of $2 \bmod t$. It follows that $r / 2$ divides both $u$ and $v$, which is impossible as $r>2$.
Note also that the factorisation

$$
2^{u v}+1=\left(2^{u}+1\right)\left(2^{u(v-1)}-2^{u(v-2)}+\ldots+2^{2 u}-2^{u}+1\right)
$$

shows that $2^{u v}+1$ is divisible by $2^{u}+1$ and $2^{v}+1$, and so is also divisible by $\left(2^{u}+1\right)\left(2^{v}+1\right) / 3$.
Let us now prove the desired result by induction on $n$. It is certainly true when $n=1$ (for example, because $2^{p_{1}}+1$ is a multiple of 3 and is at least 27 ), so we assume that $2^{p_{1} \ldots p_{n-1}}+1$ has at least $4^{n-1}$ divisors and consider $2^{p_{1} \ldots p_{n}}+1$. Setting $u=p_{1} \ldots p_{n-1}$ and $v=p_{n}$ in the above, we see that $2^{u}+1$ and $\left(2^{v}+1\right) / 3$ are coprime, whence $m=\left(2^{u}+1\right)\left(2^{v}+1\right) / 3$ has at least $2 \times 4^{n-1}$ divisors.
Now, we know that $m$ divides $2^{u v}+1$. Moreover, from $u v>2(u+v)$ when $u, v \geqslant 5$, we see that $2^{u v}+1>m^{2}$. By the fact mentioned in the comment above, it follows that $d\left(2^{u v}+1\right) \geqslant 2 d(m) \geqslant 4^{n}$, as required.

## Further comment

2. From a more advanced point of view, $f\left(p_{1} p_{2} \ldots p_{n}\right)$ is the product of cyclotomic polynomials at 2 , that is the product of $\Phi_{2 m}(2)$ over $m \mid p_{1} \ldots p_{n}$. It turns out that $\Phi_{r}(2)$ and $\Phi_{s}(2)$ are coprime unless $r / s$ is a prime power (this is not an easy fact), from which it follows that $f\left(p_{1} p_{2} \ldots p_{n}\right)$ has at least $2^{n-1}$ prime divisors. Hence $d\left(f\left(p_{1} p_{2} \ldots p_{n}\right)\right) \geqslant 2^{2^{n-1}}$, which is much more than $4^{n}$ when $n$ is large.

N4. Is there a positive integer $m$ such that the equation

$$
\frac{1}{a}+\frac{1}{b}+\frac{1}{c}+\frac{1}{a b c}=\frac{m}{a+b+c}
$$

has infinitely many solutions in positive integers $a, b, c$ ?

Solution. If $a=b=c=1$ then $m=12$, and we proceed to show that, for this fixed value of $m$, there are infinitely many solutions in positive integers $a, b, c$. Write

$$
\frac{1}{a}+\frac{1}{b}+\frac{1}{c}+\frac{1}{a b c}-\frac{12}{a+b+c}=\frac{p(a, b, c)}{a b c(a+b+c)},
$$

where $p(a, b, c)=a^{2}(b+c)+b^{2}(c+a)+c^{2}(a+b)+a+b+c-9 a b c$. Suppose that $(x, a, b)$ is a solution with $x<a<b$, that is $p(x, a, b)=0$. Then, regarding this as a quadratic equation in $x$, we see that $y=(a b+1) / x>b$ is also a solution, except that we need to establish that such a value $y$ is integral.
Let $a_{0}=a_{1}=a_{2}=1$, and define

$$
a_{n+2}=\frac{a_{n} a_{n+1}+1}{a_{n-1}}, \quad \text { for each } n \geqslant 1 .
$$

We now prove the following assertions simultaneously by induction:
(i) $a_{n-1} \mid a_{n} a_{n+1}+1$,
(ii) $a_{n} \mid a_{n-1}+a_{n+1}$,
(iii) $a_{n+1} \mid a_{n-1} a_{n}+1$.

The three assertions are true when $n=1$ from the initial values for $a_{0}, a_{1}, a_{2}$, and we suppose that they are true when $n=k$. Thus (i) implies that $a_{k-1}$ and $a_{k}$ are coprime and that $a_{k-1}$ divides $\left(a_{k} a_{k+1}+1\right) a_{k+1}+a_{k-1}$, whereas (ii) gives $a_{k} \mid a_{k} a_{k+1}^{2}+a_{k+1}+a_{k-1}$, so that together $a_{k} a_{k-1} \mid a_{k} a_{k+1}^{2}+a_{k+1}+a_{k-1}$, that is $a_{k} \mid a_{k+1}\left(a_{k} a_{k+1}+1\right) / a_{k-1}+1=a_{k+1} a_{k+2}+1$, which is (i) when $n=k+1$.

Similarly (i) also implies that $a_{k-1}$ and $a_{k+1}$ are coprime, and that $a_{k-1} \mid a_{k} a_{k+1}+1+a_{k} a_{k-1}$, whereas (iii) gives $a_{k+1} \mid a_{k} a_{k-1}+1+a_{k} a_{k+1}$, so that together $a_{k-1} a_{k+1} \mid a_{k}\left(a_{k-1}+a_{k+1}\right)+1$, that is $a_{k+1} \mid a_{k}+\left(a_{k} a_{k+1}+1\right) / a_{k-1}=a_{k}+a_{k+2}$, which is (ii) when $n=k+1$.
Finally, the definition of $a_{k+2}$ together with (i) implies $a_{k+2} \mid a_{k} a_{k+1}+1$, which is (iii) when $n=k+1$.

Therefore $\left(a_{n}\right)$ is a sequence of integers, strictly increasing from $n \geqslant 2$, and $p\left(a_{n}, a_{n+1}, a_{n+2}\right)=0$ for all $n$. In other words, $\left(a_{n}, a_{n+1}, a_{n+2}\right)$ is a solution to the given equation, with

$$
\left(a_{n}\right)=(1,1,1,2,3,7,11,26,41,97,154, \ldots)
$$

## Comments

1. Another method is to define $\left(c_{n}\right)$ by $c_{0}=2, c_{1}=3$, and $c_{2 n}=3 c_{2 n-1}-c_{2 n-2}$ and $c_{2 n+1}=2 c_{2 n}-c_{2 n-1}$, and use induction to show that the triples $\left(c_{n}, c_{n+1}, c_{n+2}\right)$ are solutions.
2. One may also apply Pell's equation to show that there are infinitely many solutions for $m=12$. Indeed, let $p(a, b, c)$ be as above. With an eye on eliminating a variable in $p$ by a substitution of the form $a+c=r b$ with a suitable $r$, we find that $p(1,1, r-1)=2(r-2)(r-3)$, showing that $r=2,3$ are suitable candidates. We therefore consider

$$
p(a, b, 2 b-a)=3 b\left(3 a^{2}-6 a b+2 b^{2}+1\right)=3 b\left(3(a-b)^{2}-b^{2}+1\right)
$$

and recall the well-known result that there are infinitely many solutions to the Pell equation $x^{2}=3 y^{2}+1$. Thus there are infinitely many positive integers $a<b$ satisfying $p(a, b, 2 b-a)=0$.
3. In fact, using a little more theory on quadratic forms, it can be shown that if the equation is soluble for a given value of $m$ then there are infinitely many solutions for that value of $m$.
4. There is nothing special about $m=12$ : there are infinitely many possible values of $m$. Indeed, the given equation may be rewritten as $m=(a+b+c)(1+a b+b c+c a) / a b c$, which becomes $m=(1+b+c)+(1+b+c)^{2} / b c$ on setting $a=1$. One can define a sequence $\left(b_{n}\right)$ with the property that $b_{n} b_{n+1}$ divides $\left(1+b_{n}+b_{n+1}\right)^{2}$; take, for example, $b_{1}=4, b_{2}=5$, set $b_{n+2}=3 b_{n+1}-b_{n}-2$, and induction then shows that $\left(b_{n+1}+b_{n}+1\right)^{2}=5 b_{n} b_{n+1}$. The corresponding value for $m$ is then $b_{n}+b_{n+1}+6$. We have one solution for this value of $m$, so by the remark above there are infinitely many solutions for this value of $m$.

N5. Let $m, n \geqslant 2$ be positive integers, and let $a_{1}, a_{2}, \ldots, a_{n}$ be integers, none of which is a multiple of $m^{n-1}$. Show that there exist integers $e_{1}, e_{2}, \ldots, e_{n}$, not all zero, with $\left|e_{i}\right|<m$ for all $i$, such that $e_{1} a_{1}+e_{2} a_{2}+\ldots+e_{n} a_{n}$ is a multiple of $m^{n}$.

Solution. Write $N$ for $m^{n}$. Let $B$ be the set of all $n$-tuples $b=\left(b_{1}, b_{2}, \ldots, b_{n}\right)$, where each $b_{i}$ is an integer with $0 \leqslant i<m$. For $b \in B$, write $f(b)$ for $b_{1} a_{1}+b_{2} a_{2}+\ldots \quad+b_{n} a_{n}$. If some distinct $b, b^{\prime} \in B$ have $f(b) \equiv f\left(b^{\prime}\right)(\bmod N)$ then we are done: setting $e_{i}=b_{i}-b_{i}^{\prime}$ we have $e_{1} a_{1}+\ldots+e_{n} a_{n} \equiv 0(\bmod N)$. So we are done unless no two $f(b)$ are congruent $\bmod N$. Since $|B|=N$, this implies that, $\bmod N$, the numbers $f(b)$ for $b \in B$ are precisely the numbers $0,1, \ldots, N-1$ (in some order). We wish to show that this is impossible.
Consider the polynomial $\sum_{b \in B} X^{f(b)}$. On the one hand, it factorises as

$$
\prod_{i=1}^{n}\left(1+X^{a_{i}}+X^{2 a_{i}}+\ldots+X^{(m-1) a_{i}}\right)
$$

but on the other hand it is equal to $1+X+X^{2}+\ldots+X^{N-1}$ whenever $X^{N}=1$. But now set $X=\exp (2 \pi i / N)$, a primitive $N$-th root of unity. Then

$$
1+X+X^{2}+\ldots+X^{N-1}=\frac{1-X^{N}}{1-X}=0
$$

but for each $i$ we have

$$
1+X^{a_{i}}+X^{2 a_{i}}+\ldots+X^{(m-1) a_{i}}=\frac{1-X^{m a_{i}}}{1-X}
$$

which is non-zero because $m a_{i}$ is not a multiple of $N$. This is a contradiction.

## Comments

1. The proof begins with a standard pigeonhole argument. The exceptional case (with each congruence class $\bmod N$ hit exactly once) is quickly identified, and looks at first glance at though it should be easily attackable. However, it is actually rather challenging. The use of the polynomial and $N$-th roots of unity is probably the most natural approach. We do not know of any bare-hands or essentially different proof.
2. The condition that no $a_{i}$ is a multiple of $m^{n-1}$ cannot be removed, as may be seen by taking $a_{i}=m^{i-1}$ for each $i$.

N6. Find all pairs of positive integers $m, n \geqslant 3$ for which there exist infinitely many positive integers $a$ such that

$$
\frac{a^{m}+a-1}{a^{n}+a^{2}-1}
$$

is itself an integer.

Solution. Suppose $m, n$ is such a pair. Clearly $n<m$.
Step 1. We claim that $f(x)=x^{m}+x-1$ is exactly divisible by $g(x)=x^{n}+x^{2}-1$ in $\mathbf{Z}[x]$. Indeed, since $g(x)$ is monic, the division algorithm gives

$$
f(x) / g(x)=q(x)+r(x) / g(x)
$$

where $\operatorname{deg}(r)<\operatorname{deg}(g)$. The remainder term $r(x) / g(x)$ tends to zero as $x \rightarrow \infty$; on the other hand it is an integer at infinitely many integers $a$. Thus $r(a) / g(a)=0$ infinitely often, and so $r \equiv 0$. The claim follows; and in particular, we note that $f(a) / g(a)$ is an integer for all integers $a$.

Step 2. Both $f(x)$ and $g(x)$ have a unique root in the interval $(0,1)$, since both functions are increasing in $[0,1]$ and span the range $[-1,1]$. Moreover it is the same root, since $g$ divides $f$; call it $\alpha$.

Step 3. We can use $\alpha$ to show that $m<2 n$. Certainly $\alpha>\phi$, where $\phi=0.618 \ldots$ is the positive root of $h(x)=x^{2}+x-1$. This is because $f$ is increasing in $(0,1)$ and $f(\phi)<h(\phi)=0=f(\alpha)$. On the other hand, if $m \geqslant 2 n$ then $1-\alpha=\alpha^{m} \leqslant\left(\alpha^{n}\right)^{2}=\left(1-\alpha^{2}\right)^{2}$, and the outer terms rearrange to give $\alpha(\alpha-1)\left(\alpha^{2}+\alpha-1\right) \geqslant 0$, which requires $\alpha \leqslant \phi$, a contradiction.

Step 4. We show that the only solution with $m<2 n$ is $(m, n)=(5,3)$. This is pure number theory, at last. Suppose we have a solution. We consider the value $a=2$, and write $d=g(2)=2^{n}+3$, so that $-2^{m} \equiv 1(\bmod d)$. Let $m=n+k$ where $1 \leqslant k<n$, so that

$$
-2^{m} \equiv\left(d-2^{n}\right) 2^{k} \equiv 3 \times 2^{k}(\bmod d),
$$

which shows that $-2^{m} \not \equiv 1(\bmod d)$ when $1 \leqslant k \leqslant n-2$. When $k=n-1$, that is $m=2 n-1$, the least positive residue $(\bmod d)$ for $-2^{m}$ is given by $3 \times 2^{n-1}-d=2^{n-1}-3$, which takes the value 1 only when $n=3$, giving $m=5$. Finally, the identity $a^{5}+a-1=\left(a^{3}+a^{2}-1\right)\left(a^{2}-a+1\right)$ shows that $(m, n)=(5,3)$ is indeed a solution.

## Comment

1. Although the above solution is entirely elementary, several separate good ideas seem to be needed to crack the problem. Step 1 is the natural way to begin, and Step 4 has several variations. Perhaps the most important-and most difficult-idea is the use of the common root $\alpha$ (in Steps 2 and 3) to obtain the quantitative bound $m<2 n$. All solutions we have seen make use of this idea in some form.

G1. Let $B$ be a point on a circle $S_{1}$, and let $A$ be a point distinct from $B$ on the tangent at $B$ to $S_{1}$. Let $C$ be a point not on $S_{1}$ such that the line segment $A C$ meets $S_{1}$ at two distinct points. Let $S_{2}$ be the circle touching $A C$ at $C$ and touching $S_{1}$ at a point $D$ on the opposite side of $A C$ from $B$. Prove that the circumcentre of triangle $B C D$ lies on the circumcircle of triangle $A B C$.

## Comments

1. In both solutions that follow, the key idea is to work with the perpendicular bisectors of $B D$ and $C D$.
2. There does not appear to be a straightforward coordinate solution.


Solution 1. Let $E$ and $F$ be the midpoints of $B D$ and $C D$ respectively, $K$ be the circumcentre of triangle $B C D$ and let $T D T^{\prime}$ be the common tangent to the two circles. Then $E K$ is perpendicular to $B D$ and bisects the angles between the tangents $B A, D T$ to $S_{1}$ at $B, D$. Hence $K$ is equidistant from $B A$ and $D T$. Similarly, $K F$ is perpendicular to $C D$ and $K$ is equidistant from $A C$ and $D T$. Hence $K$ is the centre of a circle touching $B A, A C$ and $D T$. Accordingly, $A K$ is a bisector of $\angle B A C$. But $K$ is also on the perpendicular bisector of $B C$ and it is known that this line meets the bisectors of $\angle B A C$ on the circumcircle of $A B C$.

Solution 2. We use the same notation as in the first solution.
Since the tangents at the ends of a chord are equally inclined to that chord, we have $\angle T D B=\angle A B D$ and $\angle T^{\prime} D C=\angle D C A$. Hence

$$
\begin{aligned}
\angle B D C & =180^{\circ}-\angle A B D+\angle D C A \\
& =180^{\circ}-(\angle A B C-\angle D B C)+(\angle D C B-\angle A C B) \\
& =\left(180^{\circ}-\angle A B C-\angle A C B\right)+(\angle D B C+\angle D C B) \\
& =\angle B A C+180^{\circ}-\angle B D C .
\end{aligned}
$$

Thus

$$
2 \angle B D C=180^{\circ}+\angle B A C
$$

Finally

$$
\begin{aligned}
\angle B K C & =\angle B K D+\angle D K C \\
& =2(\angle E K D+\angle D K F)=2 \angle E K F \\
& =2\left(180^{\circ}-\angle B D C\right)=180^{\circ}-\angle B A C,
\end{aligned}
$$

so that $K$ lies on circle $A B C$.

G2. Let $A B C$ be a triangle for which there exists an interior point $F$ such that $\angle A F B=\angle B F C=\angle C F A$. Let the lines $B F$ and $C F$ meet the sides $A C$ and $A B$ at $D$ and $E$ respectively. Prove that

$$
A B+A C \geqslant 4 D E .
$$

## Comments

1. We present two solutions, a geometrical one and an algebraic one, both of which use standard procedures and are of moderate difficulty.
2. Though the geometrical solution uses known properties of the Fermat point, these are very easy to deduce directly.
3. A complex variable solution is also possible because of the $120^{\circ}$ angles, but it is comparable with the other methods in length and difficulty.
4. Ptolemy's inequality applied to the quadrilateral $A D F E$ does not seem to produce the required result.

Solution 1. We need the following lemma:
Lemma. A triangle $D E F$ is given. Points $P$ and $Q$ lie on $F D, F E$ respectively, so that $P F \geqslant \lambda D F$ and $Q F \geqslant \lambda E F$, where $\lambda>0$. If $\angle P F Q \geqslant 90^{\circ}$ then $P Q \geqslant \lambda D E$.
Proof: Let $\angle P F Q=\theta$. Since $\theta \geqslant 90^{\circ}$, we have $-\cos \theta \geqslant 0$. Now, by the cosine law, we have $P Q^{2}=P F^{2}+Q F^{2}-2 \cos \theta(P F)(Q F) \geqslant(\lambda D F)^{2}+(\lambda E F)^{2}-2 \cos \theta(\lambda D F)(\lambda E F)=(\lambda D E)^{2}$ from which $P Q \geqslant \lambda D E$, as required.


We now start the main proof. Note that $\angle A F E=\angle B F E=\angle C F D=\angle A F D=60^{\circ}$. Now let the lines $B F, C F$ meet the circumcircles of triangles $C F A, A F B$ at the points $P, Q$ respectively. Then it is easy to see that both triangles $C P A$ and $A Q B$ are equilateral. We now use the lemma with $\lambda=4$ and $\theta=120^{\circ}$. To see how, let $P_{1}$ be the foot of the perpendicular from $F$ to the line $A C$ and suppose the perpendicular bisector of $A C$ meets the circumcircle $C F A$ at $P$ and $P_{2}$. Let $M$ be the midpoint of $A C$. Then $P D / D F=P M / F P_{1} \geqslant P M / M P_{2}=3$ so $P F \geqslant 4 D F$. Similarly we have $Q F \geqslant 4 E F$. Since $\angle D F E=120^{\circ}$, the lemma applies and so $P Q \geqslant 4 D E$. Finally, using the triangle inequality, $A B+A C=A Q+A P \geqslant P Q \geqslant 4 D E$.

## Further comment

5. An alternative argument may be used to prove $P F \geqslant 4 D F$. Since the area $[C F A]=[A F D]+[C F D]$ we have $(C F)(A F)=(C F)(D F)+(A F)(D F)$, from which

$$
D F=\frac{(C F)(A F)}{(C F)+(A F)}
$$

But it is easily shown, by Ptolemy's theorem for the cyclic quadrilateral $A F C P$ for example, that $C F+A F=P F$, so $P F / D F=\{(C F)+(A F)\}^{2} /\{(C F)(A F)\} \geqslant 4$.

Solution 2. Let $x, y, z$ denote the lengths of $A F, B F, C F$ respectively. Then, from (*), we have $D F=x z /(x+z)$ and similarly $E F=x y /(x+y)$. Applying the cosine law to triangles $A B F$, $A C F, D E F$ the given inequality becomes

$$
\sqrt{x^{2}+x y+y^{2}}+\sqrt{x^{2}+x z+z^{2}} \geqslant 4 \sqrt{\left(\frac{x y}{x+y}\right)^{2}+\left(\frac{x z}{x+z}\right)^{2}+\left(\frac{x y}{x+y}\right)\left(\frac{x z}{x+z}\right)}
$$

Since $(x+y) / 4 \geqslant x y /(x+y)$ and $(x+z) / 4 \geqslant x z /(x+z)$ it is sufficient to prove

$$
\sqrt{x^{2}+x y+y^{2}}+\sqrt{x^{2}+x z+z^{2}} \geqslant \sqrt{(x+y)^{2}+(x+z)^{2}+(x+y)(x+z)} .
$$

It is easy to check that the square of the left-hand side minus the square of the right-hand side comes to

$$
2 \sqrt{\left(x^{2}+x y+y^{2}\right)\left(x^{2}+x z+z^{2}\right)}-\left(x^{2}+2(y+z) x+y z\right) .
$$

It is sufficient, therefore to show that the square of the first term is greater than or equal to the square of the second term. But a short calculation shows that the difference between these two squares is equal to $3\left(x^{2}-y z\right)^{2} \geqslant 0$.

## Further comment

6. It is easy to show that equality holds if and only if triangle $A B C$ is equilateral, but there seems no interest in making this part of the question.

G3. The circle $S$ has centre $O$, and $B C$ is a diameter of $S$. Let $A$ be a point of $S$ such that $\angle A O B<120^{\circ}$. Let $D$ be the midpoint of the arc $A B$ which does not contain $C$. The line through $O$ parallel to $D A$ meets the line $A C$ at $I$. The perpendicular bisector of $O A$ meets $S$ at $E$ and at $F$. Prove that $I$ is the incentre of the triangle $C E F$.

## Comments

1. The condition $\angle A O B<120^{\circ}$ ensures that $I$ is internal to triangle $C E F$.
2. Besides the two solutions given, other proofs using circle and triangle properties are possible; a coordinate method would appear to be lengthy.


Solution 1. $A$ is the midpoint of arc $E A F$, so $C A$ bisects $\angle E C F$. Now, since $O A=O C$, $\angle A O D=\frac{1}{2} \angle A O B=\angle O A C$ so $O D$ is parallel to $I A$ and $O D A I$ is a parallelogram. Hence $A I=O D=O E=A F$ since $O E A F$ (with diagonals bisecting each other at right angles) is a rhombus. Thus

$$
\begin{aligned}
\angle I F E & =\angle I F A-\angle E F A=\angle A I F-\angle E C A \\
& =\angle A I F-\angle I C F=\angle I F C .
\end{aligned}
$$

Therefore, $I F$ bisects angle $E F C$ and $I$ is the incentre of triangle $C E F$.

Solution 2. As in the first solution, $O D A I$ is a parallelogram. Thus both $O$ and $I$ lie on the image $S^{*}$ of the circle $S$ under the half-turn about the midpoint $M$ of $E F$. Let $I_{0}$ be the incentre of the triangle CEF. Since $A$ is the midpoint of the arc $E F$ of $S$ which does not contain $C$, both $I$ and $I_{0}$ lie on the side $C A$, which is the internal bisector of $\angle E C F$. Note that

$$
A O=O E=E A=A F=F O,
$$

implying that $A E O$ and $A F O$ are congruent equilateral triangles. It follows that $\angle E O F=120^{\circ}$. Since $I_{0}$ is the incentre and $O$ the circumcentre of $C E F$ we have

$$
\angle E I_{0} F=90^{\circ}+\frac{1}{2} \angle E C F=90^{\circ}+\frac{1}{4} \angle E O F=120^{\circ} .
$$

It follows that $I_{0}$, as well as $I$, lies on $S^{*}$. Since $S^{*}$ has a unique intersection with the side $A C$, we conclude that $I=I_{0}$.

G4. Circles $S_{1}$ and $S_{2}$ intersect at points $P$ and $Q$. Distinct points $A_{1}$ and $B_{1}$ (not at $P$ or $Q$ ) are selected on $S_{1}$. The lines $A_{1} P$ and $B_{1} P$ meet $S_{2}$ again at $A_{2}$ and $B_{2}$ respectively, and the lines $A_{1} B_{1}$ and $A_{2} B_{2}$ meet at $C$. Prove that, as $A_{1}$ and $B_{1}$ vary, the circumcentres of triangles $A_{1} A_{2} C$ all lie on one fixed circle.

## Comments

1. The solution establishes the essential fact that the circle to be identified passes through $Q$ and the centres $O_{1}, O_{2}$ of $S_{1}, S_{2}$ respectively. A solver must appreciate this before composing a solution. The motivation may arise from considering certain special or limiting cases. For example, when $A_{1} P$ is tangent to $S_{2}$ at $P$ then $A_{2}$ coincides with $P$ and $C$ coincides with $B_{1}$. The circumcircle of triangle $A_{1} A_{2} C$ is then $S_{1}$ and its circumcentre $O$ coincides with $O_{1}$. Also if $B_{1}$ is close to $Q$, so are $B_{2}$ and $C$, indicating that $Q$ lies on the circle to be identified.
2. Although the solution given is short and the problem is by no means hard, it is not as straightforward as the solution may at first sight suggest (see above comment).
3. An analytic solution is possible, but the best we could manage took three full sheets of writing!


## Solution.

Step 1. The points $A_{1}, C, A_{2}, Q$ are concyclic.
Proof: We prove this by showing that the opposite angles of the quadrilateral add up to $180^{\circ}$. We have $\angle A_{1} C A_{2}+\angle A_{1} Q A_{2}=\angle A_{1} C A_{2}+\angle A_{1} Q P+\angle P Q A_{2}=\angle B_{1} C B_{2}+\angle C B_{1} B_{2}+\angle C B_{2} B_{1}=180^{\circ}$.
Here we have made use of the circle property that the exterior angle of a cyclic quadrilateral is equal to the interior opposite angle and also that angles in the same segment are equal.

Step 2. Let $O$ be the circumcentre of triangle $A_{1} A_{2} C$. Then the points $O, O_{1}, Q, O_{2}$ are concyclic.
Proof: We again prove that opposite angles of the quadrilateral add up to $180^{\circ}$.
From Step 1 we have $O Q=O A_{1}$. Also $O_{1} Q=O_{1} A_{1}$. Hence $\angle O O_{1} Q=\frac{1}{2} \angle A_{1} O_{1} Q$ $=180^{\circ}-\angle A_{1} P Q$. Similarly $\angle O O_{2} Q=180^{\circ}-\angle A_{2} P Q$. Here we have used the property that the angle at the centre is twice the angle at the circumference and the angle properties of a cyclic quadrilateral. Hence $\angle O O_{1} Q+\angle O O_{2} Q=180^{\circ}-\angle A_{1} P Q+180^{\circ}-\angle A_{2} P Q=180^{\circ}$. Thus, the centres of the circumcircles of all possible triangles $A_{1} A_{2} C$ (and similarly for triangles $\left.B_{1} B_{2} C\right)$ lie on a fixed circle through $O_{1}, O_{2}$ and $Q$.

## Further comment

4. There are some additional features about this configuration which may arise in alternative proofs. For example, if the tangents at $A_{1}, A_{2}$ meet at $C^{\prime}$ then $A_{1}, A_{2}, C, C^{\prime}$ are concyclic. Since it is easy to prove that $C^{\prime}, A_{1}, A_{2}, Q$ are concyclic, we have an alternative proof of Step 1.

G5. For any set $S$ of five points in the plane, no three of which are collinear, let $M(S)$ and $m(S)$ denote the greatest and smallest areas, respectively, of triangles determined by three points from $S$. What is the minimum possible value of $M(S) / m(S)$ ?

## Solution.

When the five points are arranged at the vertices of a regular pentagon, it is easy to check that $M(S) / m(S)$ equals the golden ratio, $\tau=(1+\sqrt{5}) / 2$. We claim that this is best possible.
Let $S$ be an arbitrary configuration, and label the points $A, B, C, D$ and $E$, so that $\triangle A B C$ has maximal area $M(S)$. In the following five steps, we prove the claim by showing that some triangle has area $M(S) / \tau$ or smaller.
Step 1. Construct a larger triangle $\triangle A^{\prime} B^{\prime} C^{\prime}$ with parallel sides to $\triangle A B C$ so that $A, B$ and $C$ lie at the midpoints of the edges $B^{\prime} C^{\prime}, C^{\prime} A^{\prime}$ and $A^{\prime} B^{\prime}$, respectively. The point $D$ must then lie on the same side of $B^{\prime} C^{\prime}$ as $B C$ otherwise $\triangle D B C$ would have greater area than $\triangle A B C$. Arguing similarly with the other edges and with the vertex $E$, it follows that both $D$ and $E$ necessarily lie within $\Delta A^{\prime} B^{\prime} C^{\prime}$ (perhaps on its boundary).


Step 2. We can assume more. Of the three triangles $\triangle A^{\prime} B C, \triangle A B^{\prime} C$ and $\triangle A B C^{\prime}$ at least one of them contains neither $D$ nor $E$. Rearranging the labels $A, B$ and $C$ if necessary, we can assume that $D$ and $E$ are contained inside the quadrilateral $B C B^{\prime} C^{\prime}$.

Step 3. Note that if an affine linear transformation of the plane is applied to the configuration $S$, the ratio $M(S) / m(S)$ remains unchanged (since all areas change by the same factor). We can therefore make the convenient assumption that $A, B$ and $C$ are vertices of a regular pentagon $A P B C Q$; if this is not already true, then a suitable affine linear transformation can be found carrying $A, B$ and $C$ to the required positions. Since $\angle A B P=\angle B A C=36^{\circ}$, it follows that $P$ lies on $B C^{\prime}$. Similarly, $Q$ lies on $C B^{\prime}$.


Step 4. If either $D$ or $E$ lies in the pentagon $A P B C Q$, then we are done. We argue for $D$ as follows: Note that $\triangle A P B$ has area $M(S) / \tau$. If $D$ lies in $\triangle A P B$, then $\triangle D P B$ has area at most $M(S) / \tau$. Likewise we are done if $D$ lies in $\triangle A Q C$. Finally if $D$ is contained in $\triangle A B C$, then one of $\triangle D A B, \triangle D B C$ or $\triangle D C A$ has area at most $M(S) / 3<M(S) / \tau$. Similarly for $E$.

Step 5. What remains is the case where $D$ and $E$ are contained in the union of the triangles $\triangle A P C^{\prime}$ and $\triangle A Q B^{\prime}$. Then $|A E|,|A D| \leqslant|A P|=|A Q|$, and on the other hand the angle $\theta=\angle E A D$ satisfies one of $0<\theta \leqslant 36^{\circ}$ (if $E$ and $D$ lie in the same triangle) or $108^{\circ} \leqslant \theta<180^{\circ}$ (if they lie in different triangles). Either way, we have Area $(\triangle A D E)=\frac{1}{2}|A D||A E| \sin \theta \leqslant \frac{1}{2}|A P||A Q| \sin 108^{\circ}=\operatorname{Area}(\triangle A P Q)=M(S) / \tau$. This completes the proof that the minimum value of $M(S) / m(S)$ is $\tau$.

## Comments

1. The difficulty is in knowing where to begin. The winning configuration (a regular pentagon) is certainly eminently guessable, but what next? It is natural to look at a largest (or smallest) triangle and work from there. After that, naive case-checking or diagram-chasing doesn't seem to work very well. The crucial observation is in Step 3, when we note that $\triangle A B C$ can be identified with part of a regular pentagon. Now the case-checking and diagram-chasing becomes comparatively clean, since the known geometry of the pentagon can be used as a reference.
2. Without something like Step 3 the problem is forbidding. It is still possible, but quite difficult, to find a clean argument - most attempts are likely to be messy and/or incomplete.
3. Reading the above proof carefully, it is easy to show that the minimum is attained precisely when $S$ is an affine linear transformation of the vertices of a regular pentagon.
4. The proof above in no way generalises when the number $n$ of points is greater than 5. It would be extremely interesting if a contestant were to find a proof that did work for some other values of $n$. For general $n$, the answer is unknown, and not even known asymptotically; this is related to the famous Heilbronn problem on the smallest triangle formed from $n$ points in the unit square.

G6. Let $n \geqslant 3$ be a positive integer. Let $C_{1}, C_{2}, C_{3}, \ldots, C_{n}$ be unit circles in the plane, with centres $O_{1}, O_{2}, O_{3}, \ldots, O_{n}$ respectively. If no line meets more than two of the circles, prove that

$$
\sum_{1 \leqslant i<j \leqslant n} \frac{1}{O_{i} O_{j}} \leqslant \frac{(n-1) \pi}{4} .
$$

## Comments

1. We present a solution, which, though fairly short, requires considerable ingenuity to devise. The question seems medium to hard in difficulty.
2. The last part of the solution is a double-counting argument, and doubtless there are many equivalent formulations possible.

Solution. We use the following Lemma.
Lemma. Let $\Omega$ be a circle of radius $\rho$ and $P R, Q S$ two chords intersecting at $X$, so that $\angle P X Q=\angle R X S=2 \alpha$. Then $\operatorname{arc} P Q+\operatorname{arc} R S=4 \alpha \rho$. (See Diagram 1.)


Diagram 1
Proof: Let $O$ be the centre of $\Omega$. Let $\angle P O Q=2 \lambda$ and $\angle R O S=2 \mu$; then $\angle Q S P=\lambda$ and $\angle R P S=\mu$, since the angle at the centre is twice the angle at the circumference. Then $\angle R X S=2 \alpha=\lambda+\mu$ and $\operatorname{arc} P Q+\operatorname{arc} R S=2 \lambda \rho+2 \mu \rho=4 \alpha \rho$.

We now start the main proof.
Surround all the given circles with a large circle $\Omega$ of radius $\rho$. Consider two circles $C_{i}, C_{j}$, with centres $O_{i}, O_{j}$ respectively. From the given condition $C_{i}$ and $C_{j}$ do not intersect. Let $2 \alpha$ be the angle between their two internal common tangents $P R, Q S$ (see Diagram 2). We have $O_{i} O_{j}=2 \operatorname{cosec} \alpha$, so $\alpha \geqslant \sin \alpha=2 / O_{i} O_{j}$.


Diagram 2
Now, from the lemma, arc $P Q+\operatorname{arc} R S=4 \alpha \rho \geqslant 8 \rho O_{i} O_{j}$, so that

$$
\frac{1}{O_{i} O_{j}} \leqslant \frac{\operatorname{arc} P Q+\operatorname{arc} R S}{8 \rho} .
$$

We now wish to consider the sum of all these arc lengths as $i, j$ range over all pairs, and we claim that any point of $\Omega$ is covered by such arcs at most $(n-1)$ times. To see this, let $T$ be any point of $\Omega$ and $T U$ a half-line tangent to $\Omega$, as in Diagram 3. Consider this half-line as it is rotated about $T$ as shown. At some stage it will intersect a pair of circles for the first time. Relabel these circles $C_{1}$ and $C_{2}$. The half-line can never intersect three circles, so at some further stage intersection with one of these circles, say $C_{1}$, is lost and the half-line will never meet $C_{1}$ again during its transit. Continuing in this way and relabelling the circles conveniently, the maximum number of times the half-line can intersect pairs of circles is ( $n-1$ ), namely when it intersects $C_{1}$ and $C_{2}, C_{2}$ and $C_{3}, \ldots, C_{n-1}$ and $C_{n}$. As $T$ was arbitrary, it follows that the sum of all the arc lengths is less than or equal to $2(n-1) \pi \rho$, and hence

$$
\sum_{1 \leqslant i<j \leqslant n} \frac{1}{O_{i} O_{j}} \leqslant \frac{(n-1) \pi}{4} .
$$



Diagram 3

## Further comment

3. If the lemma proves elusive, a solver could construct a proof in which $\Omega$ is sufficiently large for the intersection points to be close to its centre, thus removing any need for the lemma.

G7. The incircle $\Omega$ of the acute-angled triangle $A B C$ is tangent to $B C$ at $K$. Let $A D$ be an altitude of triangle $A B C$ and let $M$ be the midpoint of $A D$. If $N$ is the other common point of $\Omega$ and $K M$, prove that $\Omega$ and the circumcircle of triangle $B C N$ are tangent at $N$.

## Comments

1. We give two solutions, both of which involve a mixture of pure geometry and computation. The problem is difficult, but not excessively so.
2. In the first solution, the point $P$ is defined as the point of intersection of $N K$ and the perpendicular bisector of $B C$, and is shown to lie on the circumcircle of triangle $B C N$ by proving $(N K)(K P)=(B K)(K C)$.
3. In the second solution, the point $P$ is defined as the intersection (other than $N$ ) of $N K$ and the circumcircle of triangle $B C N$, and is shown to lie on the perpendicular bisector of $B C$ by proving that $N K$ bisects $\angle B N C$.
4. In the two solutions, we perform some manipulations that only make sense when $A B$ is not equal to $A C$. This is why we start by dealing with the (trivial) case when $A B=A C$. It would be possible to add the words 'non-isosceles' in the statement of the problem, but we feel that this would detract from its elegance, especially as the result does still hold in the isosceles case.


Solution 1. We may assume that $A B \neq A C$, as if $A B=A C$ then the result is trivial (as the distance between the centres of the two internally tangent circles is equal to the difference of their radii). By symmetry, we may assume that $A B<A C$.
Let the perpendicular bisector of the side $B C$ intersect $N K$ and $B C$ at $P$ and $A^{\prime}$ respectively. It is sufficient to prove that $N$, the incentre $I$ of triangle $A B C$ and $S$, the circumcentre of triangle $B C N$, are collinear. Since $I K$ and $S P$ are parallel, both being perpendicular to $B C$, it is sufficient to prove that $P$ lies on the circumcircle of triangle $B C N$; for once we know $S P=S N$ then $\angle P N S=\angle N P S=\angle N K I=\angle P N I$, and NIS is a straight line. To establish what is wanted we show $(N K)(K P)=(B K)(K C)$.

Using the standard notation for triangle $A B C$ (with $s=(a+b+c) / 2$ ), we have $B K=s-b$ and $K C=s-c$, so $(B K)(K C)=(s-b)(s-c)$. By the cosine law for triangle $A B C$, we have $\cos B=\left(c^{2}+a^{2}-b^{2}\right) / 2 c a$ and then $B D=c \cos B=\left(c^{2}+a^{2}-b^{2}\right) / 2 a$. Now $K A^{\prime}=B A^{\prime}-B K=\frac{1}{2}(b-c)$ and $D K=B K-B D=(b-c)(s-a) / a$. Let $\angle M K D=\phi$; then

$$
\tan \phi=\frac{M D}{D K}=\frac{\frac{1}{2}(A D) a}{(b-c)(s-a)}=\frac{[A B C]}{(b-c)(s-a)},
$$

where [ $A B C$ ] is the area of $A B C$. Now $\angle N I K=2 \phi$, so $N K=2 r \sin \phi$, where $r$ is the inradius of $A B C$. Finally, from triangle $A^{\prime} K P$ we have $K P=K A^{\prime} \sec \phi$ and hence

$$
(N K)(K P)=2 r\left(K A^{\prime}\right) \tan \phi=\frac{r[A B C]}{(s-a)}=\frac{[A B C]^{2}}{s(s-a)}=(s-b)(s-c)=(B K)(K C) .
$$

Here we have used the well-known expressions for area: $[A B C]=r s=\sqrt{s(s-a)(s-b)(s-c)}$.

Solution 2. As in Solution 1, we may assume $A B<A C$, and it is sufficient to show that NIS is a straight line. But now we define $P$ to be the intersection (other than $N$ ) of $N K$ with the circumcircle of triangle $B N C$. Now $S P=S N$ implies $\angle S P N=\angle S N P$ and $I N=I K$ implies $\angle I K N=\angle I N K$, and NIS is a straight line if and only if all these angles are equal, which is when $I K$ and $S P$ are parallel. Since $I K$ is perpendicular to $B C$ this means that $S P$ must be also, and hence it is sufficient to show that $P$ is the midpoint of the $\operatorname{arc} B C$. To establish this, we show that $N K P$ bisects $\angle B N C$ for which it is sufficient to show that $B N / C N=B K / C K$. Again let $\angle M K D=\phi$. Now, by the cosine rule,

$$
B N^{2}=N K^{2}+B K^{2}-2(N K)(B K) \cos \phi
$$

and

$$
C N^{2}=N K^{2}+C K^{2}+2(N K)(C K) \cos \phi .
$$

So it is sufficient to show

$$
\frac{B K^{2}}{C K^{2}}=\frac{N K^{2}+B K^{2}-2(N K)(B K) \cos \phi}{N K^{2}+C K^{2}+2(N K)(C K) \cos \phi}
$$

or $(C K-B K) N K=2(B K)(C K) \cos \phi$.
Now $N K=2 r \sin \phi$, so it is sufficient to prove $2 r(C K-B K) \tan \phi=2(B K)(C K)$. But $\tan \phi=M D / D K=\frac{1}{2} A D / D K=\frac{1}{2} c \sin B /(s-b-c \cos B)$, since $M$ is the midpoint of $A D$. Now $B K=r \cot \frac{1}{2} B$ and $C K=r \cot \frac{1}{2} C$, so it is sufficient to prove

$$
\left(\cot \frac{1}{2} C-\cot \frac{1}{2} B\right)(c \sin B) /(a+c-b-2 c \cos B)=\cot \frac{1}{2} B \cot \frac{1}{2} C .
$$

Using the sine rule and $a=c \cos B+b \cos C$ this reduces to proving that

$$
\begin{equation*}
\sin C \sin B\left(\cot \frac{1}{2} C-\cot \frac{1}{2} B\right)=\cot \frac{1}{2} B \cot \frac{1}{2} C(\sin C-\sin B+\sin B \cos C-\sin C \cos B) . \tag{*}
\end{equation*}
$$

Putting $\left(^{*}\right)$ into half-angles, and cancelling $\sin \frac{1}{2}(B-C)$ this resolves to

$$
\sin B \sin C=4 \sin \frac{1}{2} B \sin \frac{1}{2} C \cos \frac{1}{2} B \cos \frac{1}{2} C,
$$

which is true.

## Further comment

5. Slight changes in the text are necessary in Solution 2 when $A B>A C$, but the solution is essentially the same.

G8. Let $S_{1}$ and $S_{2}$ be circles meeting at the points $A$ and $B$. A line through $A$ meets $S_{1}$ at $C$ and $S_{2}$ at $D$. Points $M, N, K$ lie on the line segments $C D, B C, B D$ respectively, with $M N$ parallel to $B D$ and $M K$ parallel to $B C$. Let $E$ and $F$ be points on those arcs $B C$ of $S_{1}$ and $B D$ of $S_{2}$ respectively that do not contain $A$. Given that $E N$ is perpendicular to $B C$ and $F K$ is perpendicular to $B D$ prove that $\angle E M F=90^{\circ}$.

## Comments

1. In the solution, the lemma looks elaborate but merely formalizes the 'obvious' similarity of two figures involving circular arcs. This seems worth making explicit as it appears to be the key to the problem.
2. A coordinate approach would be impracticable.


## Solution.

Lemma. If $P_{1} Q_{1} R_{1}$ and $P_{2} Q_{2} R_{2}$ are circular arcs with $\angle P_{1} Q_{1} R_{1}=\angle P_{2} Q_{2} R_{2}$ and $T_{1}, T_{2}$ are the feet of the perpendiculars from $Q_{1}, Q_{2}$ to $P_{1} R_{1}, P_{2} R_{2}$ respectively, then if $P_{1} T_{1} / T_{1} R_{1}=P_{2} T_{2} / T_{2} R_{2}$ then the triangles $P_{1} Q_{1} R_{1}, P_{2} Q_{2} R_{2}$ are similar.
Proof: If $Q_{2}{ }^{\prime}$ is the unique point on $\operatorname{arc} P_{2} Q_{2} R_{2}$ making triangles $P_{1} Q_{1} R_{1}, P_{2} Q_{2}{ }^{\prime} R_{2}$ equiangular and therefore similar, and if $Q_{2}{ }^{\prime} T_{2}{ }^{\prime}$ is
 perpendicular to $P_{2} R_{2}$, then $P_{2} T_{2}^{\prime} / T_{2}^{\prime} R_{2}=P_{1} T_{1} / T_{1} R_{1}=P_{2} T_{2} / T_{2} R_{2}$, so $T_{2}{ }^{\prime}=T_{2}$ and $Q_{2}{ }^{\prime}=Q_{2}$.
Turning now to the problem we have

$$
\begin{array}{rlr}
B N / N C & =D M / M C \quad \text { since } M N \| D B \\
& =D K / K B \quad \text { since } M K \| C B .
\end{array}
$$

Let $F K$ produced meet $S_{2}$ again at $Q$. Then $\angle B Q D=\angle B A D=\angle B E C$. By the Lemma, triangles $B E C, D Q B$ are
 similar. Hence $\angle E B C=\angle Q D B=\angle Q F B$ and the right-angled triangles $B N E, F K B$ are similar.
Now $\angle M N B=\angle M K B$ since $M K B N$ is a parallelogram, so $\angle E N M=\angle M K F$. Also $\frac{M N}{K F}=\frac{B K}{K F}=\frac{E N}{N B}=\frac{E N}{M K}$. Therefore triangles $E N M, M K F$ are similar and $\angle N M E=\angle K F M$. Since lines $M N, K F$ are perpendicular, so are $E M$ and $F M$.

## Further comment

3. From the similarity of the right-angled triangles $B N E, F K B$ it follows easily that $\angle E A F=90^{\circ}$ as well.

A1. Find all functions $f$ from the reals to the reals such that

$$
f(f(x)+y)=2 x+f(f(y)-x)
$$

for all real $x, y$.

Solution. For each real $c$, the function given by $f(x)=x+c$ is a solution for the given functional equation, since it makes both sides equal $x+y+2 c$. We claim that these are the only solutions.
Our strategy is to derive an equation of the form $f(X)=X+c$, where $X$ is an expression whose values are guaranteed to run over all real numbers.
We claim first that $f$ is surjective. Indeed, set $y=-f(x)$ in the functional equation. This gives

$$
f(0)=2 x+f(f(-f(x))-x)
$$

or

$$
f(0)-2 x=f(f(-f(x))-x) .
$$

As all real numbers have the form $f(0)-2 x$, for each real $y$ there is a $z$ with $y=f(z)$, as claimed. In particular there is an $a$ with $f(a)=0$. Set $x=a$ in the functional equation. This gives

$$
f(y)=2 a+f(f(y)-a),
$$

or equivalently

$$
f(y)-a=f(f(y)-a)+a .
$$

As $f$ is surjective, for each real $x$ there is a real $y$ with $x=f(y)-a$. Hence

$$
x=f(x)+a
$$

for all $x$, that is $f(x)=x-a$.

## Comment

1. This is an easy problem. It seems to be crucial to note that $f$ is surjective.

A2. Let $a_{1}, a_{2}, \ldots$ be an infinite sequence of real numbers, for which there exists a real number $c$ with $0 \leqslant a_{i} \leqslant c$ for all $i$, such that

$$
\left|a_{i}-a_{j}\right| \geqslant \frac{1}{i+j} \text { for all } i, j \text { with } i \neq j .
$$

Prove that $c \geqslant 1$.

Solution. Fix $n \geqslant 2$, and let $\sigma(1), \sigma(2), \ldots, \sigma(n)$ be the permutation of $1,2, \ldots, n$ which orders the first $n$ elements of the sequence:

$$
0 \leqslant a_{\sigma(1)}<a_{\sigma(2)}<\ldots<a_{\sigma(n)} \leqslant c .
$$

Then

$$
\begin{align*}
c & \geqslant a_{\sigma(n)}-a_{\sigma(1)} \\
& =\left(a_{\sigma(n)}-a_{\sigma(n-1)}\right)+\left(a_{\sigma(n-1)}-a_{\sigma(n-2)}\right)+\ldots+\left(a_{\sigma(2)}-a_{\sigma(1)}\right) \\
& \geqslant \frac{1}{\sigma(n)+\sigma(n-1)}+\frac{1}{\sigma(n-1)+\sigma(n-2)}+\ldots+\frac{1}{\sigma(2)+\sigma(1)} . \tag{*}
\end{align*}
$$

Now, using the Cauchy-Schwarz inequality, we obtain

$$
\left(\frac{1}{\sigma(n)+\sigma(n-1)}+\ldots+\frac{1}{\sigma(2)+\sigma(1)}\right)((\sigma(n)+\sigma(n-1))+\ldots+(\sigma(2)+\sigma(1))) \geqslant(n-1)^{2}
$$

so

$$
\begin{aligned}
\frac{1}{\sigma(n)+\sigma(n-1)}+\ldots+\frac{1}{\sigma(2)+\sigma(1)} & \geqslant \frac{(n-1)^{2}}{2(\sigma(1)+\ldots+\sigma(n))-\sigma(1)-\sigma(n)} \\
& =\frac{(n-1)^{2}}{n(n+1)-\sigma(1)-\sigma(n)} \\
& \geqslant \frac{(n-1)^{2}}{n^{2}+n-3} \\
& \geqslant \frac{n-1}{n+3} .
\end{aligned}
$$

From (*) it follows that the inequality

$$
c \geqslant \frac{n-1}{n+3}=1-\frac{4}{n+3}
$$

holds for all $n \geqslant 2$. Thus we must have $c \geqslant 1$.

## Comments

1. What makes the question challenging is: how do we bring in the value of $c$ ? Which bits of the data should we use? The key step is to realise that to make use of $c$ all we need are the distances between adjacent terms. Having got equation (*), the rest is then easy, and there are several ways to finish off the proof.
2. We do not know what the smallest value of $c$ actually is.
3. The solution relies on finding a lower bound for the quantity

$$
\frac{1}{\sigma(1)+\sigma(2)}+\frac{1}{\sigma(2)+\sigma(3)}+\ldots+\frac{1}{\sigma(n-1)+\sigma(n)}
$$

where $\sigma$ is an arbitrary permutation of $(1,2, \ldots, n)$. An alternative would be to take this as the heart of the question, and ask for the exact minimum, thus:

A2'. What is the minimum value of

$$
\frac{1}{\sigma(1)+\sigma(2)}+\frac{1}{\sigma(2)+\sigma(3)}+\ldots+\frac{1}{\sigma(n-1)+\sigma(n)}
$$

as $\sigma$ ranges over all permutations of $\{1,2, \ldots, n\}$ ?

The optimal permutation turns out to be the one given by $\sigma(1)=1, \sigma(n)=2, \sigma(2)=3$, $\sigma(n-1)=4$, and so on. To prove this, we use induction, but it is vital to prove a stronger statement: that if we look at permutations of a general sequence $x_{1}, x_{2}, \ldots, x_{n}$ instead of just $1,2, \ldots, n$ (where say $x_{1}<x_{2}<\ldots<x_{n}$ ), then the optimal permutation is again $\sigma\left(x_{1}\right)=x_{1}, \sigma\left(x_{n}\right)=x_{2}, \sigma\left(x_{2}\right)=x_{3}, \sigma\left(x_{n-1}\right)=x_{4}$, and so on. The proof is more difficult, and more interesting, than that of A2. The only drawback is that A2' lacks the 'how on earth can we make use of the information?' puzzle that contestants face with A2.

A3. Let $P$ be a cubic polynomial given by $P(x)=a x^{3}+b x^{2}+c x+d$, where $a, b, c, d$ are integers and $a \neq 0$. Suppose that $x P(x)=y P(y)$ for infinitely many pairs $x, y$ of integers with $x \neq y$. Prove that the equation $P(x)=0$ has an integer root.

## Comment

1. The main ideas in the solution are that $x+y$ is bounded for all solutions $x, y$ (consider the shape of the quartic $x P(x))$, and that $P$ is then symmetric about one particular value of $(x+y) / 2$ which is taken infinitely often.

Solution. Let $x, y$ be distinct integers satisfying $x P(x)=y P(y)$ so that

$$
\begin{gathered}
x\left(a x^{3}+b x^{2}+c x+d\right)=y\left(a y^{3}+b y^{2}+c y+d\right) \\
\text { i.e. } a\left(x^{4}-y^{4}\right)+b\left(x^{3}-y^{3}\right)+c\left(x^{2}-y^{2}\right)+d(x-y)=0 .
\end{gathered}
$$

Dividing by $x-y(\neq 0)$ leads to

$$
\begin{equation*}
a\left(x^{3}+x^{2} y+x y^{2}+y^{3}\right)+b\left(x^{2}+x y+y^{2}\right)+c(x+y)+d=0 . \tag{1}
\end{equation*}
$$

It is convenient to write

$$
\begin{equation*}
s=x+y, \quad t=x y \tag{2}
\end{equation*}
$$

Since

$$
x^{3}+x^{2} y+x y^{2}+y^{3}=(x+y)\left(x^{2}+y^{2}\right)=s\left(s^{2}-2 t\right)
$$

and

$$
x^{2}+x y+y^{2}=s^{2}-t,
$$

(1) can be written in the form

$$
a s\left(s^{2}-2 t\right)+b\left(s^{2}-t\right)+c s+d=0
$$

or equivalently as

$$
\begin{equation*}
P(s)=(2 a s+b) t . \tag{3}
\end{equation*}
$$

We claim that the integer $s$ can take only finitely many values. Indeed, consider the right-hand side of (3). Since $s^{2}-4 t=(x-y)^{2} \geqslant 0$, we have $|t|<s^{2} / 4$ so that $|(2 a s+b) t| \leqslant\left|(2 a s+b)\left(s^{2} / 4\right)\right|$.
Equation (3) therefore leads to

$$
\left|a s^{3}+b s^{2}+c s+d\right| \leqslant\left|\frac{a}{2} s^{3}+\frac{b}{4} s^{2}\right|
$$

which can only be true for finitely many values of the integer $s$, as required.
Write $Q(x)$ for $x P(x)$. The equation $x P(x)=y P(y)$ becomes $Q(x)=Q(y)$, which holds for infinitely many pairs of distinct integers $x, y$. Equivalently, $Q(r)=Q(s-r)$ holds for infinitely many pairs of integers $s, r$ with $s$ as in (2). But $s$ can only take finitely many values. Hence for (at least) one integer $s$, the equation $Q(r)=Q(s-r)$ must hold for infinitely many integers. But then $Q(x)$ and $Q(s-x)$ are polynomials of degree 4 which are equal for infinitely many integer values of $x$. They must therefore be equal for all real numbers $x$.

To finish the proof, we consider two cases.
Case 1: $s \neq 0$. We have $x P(x)=(s-x) P(s-x)$ for all real numbers $x$. Take $x=s$ to get $s P(s)=0$ so that $P(s)=0$ as $s \neq 0$. Hence $x=s$ is an integer root of $P(x)=0$.
Case 2: $s=0$. We now have $Q(x)=Q(-x)$ for all real numbers $x$ so that $Q$ is an even function. As $Q(x)$ is divisible by $x$, it must be divisible by $x^{2}$ i.e. $Q(x)=x P(x)=x^{2} R(x)$ for some polynomial $R(x)$. Hence $P(x)=x R(x)$ so that $P(0)=0$. Again, the equation $P(x)=0$ has an integer root, namely $x=0$.

## Further comment

2. By examining further Cases 1 and 2 above, it is not too hard to show that polynomials $P$ satisfying the conditions of the problem have the general form

$$
P(x)=(x-k)\left(a x^{2}-a k x+m\right)
$$

where $a, k$ and $m$ are integers.

A4. Find all functions $f$ from the reals to the reals such that

$$
(f(x)+f(z))(f(y)+f(t))=f(x y-z t)+f(x t+y z)
$$

for all real $x, y, z, t$.

Solution. We are given that

$$
\begin{equation*}
f(x y-z t)+f(x t+y z)=(f(x)+f(z))(f(y)+f(t)) \tag{*}
\end{equation*}
$$

for all real $x, y, z, t$. The equation (*) has the solutions $f(x)=0$ for all $x, f(x)=1 / 2$ for all $x$ and $f(x)=x^{2}$ for all $x$. These make both sides of $\left({ }^{*}\right)$ equal to 0 , to 1 and to $\left(x^{2}+z^{2}\right)\left(y^{2}+t^{2}\right)$ respectively. We claim that there are no other solutions.
Suppose (*) holds. Then setting $x=y=z=0$ gives $2 f(0)=2 f(0)(f(0)+f(t))$. In particular $2 f(0)=4 f(0)^{2}$ and $\operatorname{so} f(0)=0$ or $f(0)=1 / 2$. If $f(0)=1 / 2$ we get $f(0)+f(t)=1$ and so $f$ is identically $1 / 2$.
Suppose then that $f(0)=0$. Then setting $z=t=0$ in (*) gives $f(x y)=f(x) f(y)$, that is $f$ is multiplicative. In particular $f(1)=f(1)^{2}$ and $\operatorname{so} f(1)=0$ or 1 . If $f(1)=0$ then $f(x)=f(x) f(1)=0$ for all $x$.
So we may suppose that $f(0)=0$ and $f(1)=1$. Setting $x=0$ and $y=t=1,\left({ }^{*}\right)$ gives

$$
f(-z)+f(z)=2 f(z)
$$

and $\operatorname{so} f(-z)=f(z)$ for each $z$, that is $f$ is an even function. So it suffices to show that $f(x)=x^{2}$ for positive $x$. Note that $f\left(x^{2}\right)=f(x)^{2} \geqslant 0$; as $f$ is an even function, $f(y) \geqslant 0$ for all $y$.
Now put $t=x$ and $z=y$ in (*) to get

$$
f\left(x^{2}+y^{2}\right)=(f(x)+f(y))^{2} .
$$

This shows that $f\left(x^{2}+y^{2}\right) \geqslant f(x)^{2}=f\left(x^{2}\right)$. Hence if $u \geqslant v \geqslant 0$ then $f(u) \geqslant f(v)$, that is $f$ is an increasing function on the positive reals.
Set $y=z=t=1$ in $(*)$ to yield

$$
f(x-1)+f(x+1)=2(f(x)+1) .
$$

By induction on $n$, it readily follows from this that $f(n)=n^{2}$ for all non-negative integers $n$. As $f$ is even, $f(n)=n^{2}$ for all integers $n$, and further, as $f$ is multiplicative, $f(a)=a^{2}$ for all rationals $a$. Suppose that $f(x) \neq x^{2}$ for some positive $x$. If $f(x)<x^{2}$ take a rational $a$ with $x>a>\sqrt{f(x)}$. Then $f(a)=a^{2}>f(x)$, but $f(a) \leqslant f(x)$ as $f$ is increasing. This is a contradiction. A similar argument shows that $f(x)>x^{2}$ is impossible. Thus $f(x)=x^{2}$ for all positive $x$, and since $f$ is even, $f(x)=x^{2}$ for all real $x$.

## Comments

1. This is a medium difficulty problem, requiring no really clever ideas, but a willingness to experiment with the functional equation to squeeze out diverse consequences. One also needs the passage from knowing $f$ on the rationals to knowing it on the reals: here the key point is that we know that $f$ is increasing (on the positive reals).
2. This problem is clearly inspired by the famous identity

$$
\left(x^{2}+z^{2}\right)\left(y^{2}+t^{2}\right)=(x y-z t)^{2}+(x t+y z)^{2}
$$

used to study sums of two squares in number theory.
3. In its original form, the problem had the variable $t$ set equal to 1 , thus:

A4'. Find all functions $f$ from the reals to the reals such that

$$
(f(x)+f(z))(f(y)+1)=f(x y-z)+f(x+y z)
$$

for all real $x, y, z$.
This has the 'sums of two squares' hidden, so that it may be less clear that $f(x)=x^{2}$ is a solution. We feel that this version is less elegant and attractive than the form we have given. It is quite easy to transform either problem into the other.

A5. Let $n$ be a positive integer that is not a perfect cube. Define real numbers $a, b, c$ by

$$
a=\sqrt[3]{n}, \quad b=\frac{1}{a-[a]}, \quad c=\frac{1}{b-[b]}
$$

where $[x]$ denotes the integer part of $x$.
Prove that there are infinitely many such integers $n$ with the property that there exist integers $r, s, t$, not all zero, such that $r a+s b+t c=0$.

Solution. Note first that it is sufficient to find rational numbers $r, s$, $t$ not all zero such that $r a+s b+t c=0$.
Let $m=[a]$ and $k=n-m^{3}$. Then

$$
1 \leqslant k \leqslant\left((m+1)^{3}-1\right)-m^{3}=3 m(m+1)
$$

From the factorisation $a^{3}-m^{3}=(a-m)\left(a^{2}+a m+m^{2}\right)$ we have

$$
b=\frac{1}{a-m}=\frac{a^{2}+a m+m^{2}}{k} .
$$

Since $a<m+1$, the numerator is less than

$$
(m+1)^{2}+(m+1) m+m^{2}=3 m^{2}+3 m+1
$$

To simplify the calculation, we shall assume that $[b]=1$. This is true provided that

$$
\begin{equation*}
3 m^{2}+3 m+1<2 k \tag{*}
\end{equation*}
$$

Now

$$
b-[b]=b-1=\frac{a^{2}+a m+m^{2}-k}{k} .
$$

Factorise the numerator in the form $a^{2}+a m+m^{2}-k=(a-x)(a-y)$, so that $x+y=-m$ and $x y=m^{2}-k$. Note that $x, y$ are real since the discriminant $(x-y)^{2}$ is $m^{2}-4\left(m^{2}-k\right)=4 k-3 m^{2}>0$ by $(*)$. Then

$$
c=\frac{1}{b-1}=\frac{k}{(a-x)(a-y)}=\frac{k\left(a^{2}+a x+x^{2}\right)\left(a^{2}+a y+y^{2}\right)}{\left(a^{3}-x^{3}\right)\left(a^{3}-y^{3}\right)} .
$$

(Note that $a^{2}+a x+x^{2}$ and $a^{2}+a y+y^{2}$ are strictly positive.)
Since

$$
x^{3}+y^{3}=(x+y)\left[(x+y)^{2}-3 x y\right]=-m\left[m^{2}-3\left(m^{2}-k\right)\right]=m\left(2 m^{2}-3 k\right)
$$

and

$$
x^{3} y^{3}=(x y)^{3}=\left(m^{2}-k\right)^{3}
$$

are integers, so is $l=\left(a^{3}-x^{3}\right)\left(a^{3}-y^{3}\right)=n^{2}-\left(x^{3}+y^{3}\right) n+x^{3} y^{3}$. Then

$$
\begin{aligned}
c & =\frac{k}{l}\left(a^{2}+a x+x^{2}\right)\left(a^{2}+a y+y^{2}\right) \\
& =\frac{k}{l}\left(a^{4}+(x+y) a^{3}+\left(x^{2}+x y+y^{2}\right) a^{2}+x y(x+y) a+x^{2} y^{2}\right) \\
& =\frac{k}{l}\left(k a^{2}+\left(m\left(k^{2}-m\right)+n\right) a+\left(m^{2}-k\right)^{2}-n m\right) .
\end{aligned}
$$

To make $r a+s b+t c=0$, choose $s$ and $t$ so that the coefficients of $a^{2}$ and $a^{0}$ vanish i.e.

$$
\frac{s}{k}+\frac{t k^{2}}{l}=0 \quad \text { and } \quad \frac{s m^{2}}{k}+\frac{t k}{l}\left(\left(m^{2}-k\right)^{2}-n m\right)=0 .
$$

The first equation gives $s=-\frac{t k^{3}}{l}$ and the second then becomes

$$
\begin{aligned}
& -\frac{t k^{2} m^{2}}{l}+\frac{t k}{l}\left(\left(m^{2}-k\right)^{2}-m n\right)=0 \\
& \text { i.e. } \frac{t k}{l}\left[\left(m^{2}-k\right)^{2}-m n-k m^{2}\right]=0
\end{aligned}
$$

The bracket simplifies as

$$
\begin{aligned}
m^{4}-2 k m^{2}+k^{2}-m n-k m^{2} & =m\left(m^{3}-n\right)-3 k m^{2}+k^{2} \\
& =-m k-3 k m^{2}+k^{2} \\
& =k\left(k-3 m^{2}-m\right) .
\end{aligned}
$$

Choose $k=3 m^{2}+m$, which satisfies $\left(^{*}\right)$ and also satisfies $1 \leqslant k \leqslant 3 m(m+1)$. We have therefore shown that for integers $n$ of the form $m^{3}+3 m^{2}+m$ we can obtain non-zero rational values of $s, t$ so that $s b+t c$ is a rational multiple of $a$. In view of the opening comment, this completes the proof.

## Comment

1. One of the key ideas is to force $b$ to have integer part 1 : this greatly simplifies what is to come. But there are still more ideas needed: pure brute-force calculation would be doomed to failure.

A6. Let $A$ be a non-empty set of positive integers. Suppose that there are positive integers $b_{1}, \ldots b_{n}$ and $c_{1}, \ldots, c_{n}$ such that
(i) for each $i$ the set $b_{i} A+c_{i}=\left\{b_{i} a+c_{i}: a \in A\right\}$ is a subset of $A$, and
(ii) the sets $b_{i} A+c_{i}$ and $b_{j} A+c_{j}$ are disjoint whenever $i \neq j$.

Prove that

$$
\frac{1}{b_{1}}+\ldots+\frac{1}{b_{n}} \leqslant 1 .
$$

## Comment

1. In the following proof, the key idea, after trying an example with the $b_{i}$ equal, is to weight the number of times each $b_{i}$ appears (this is the use of the $p_{i}$ below). After that, the calculations are quite easy, and there are several ways to accomplish them. However, this key idea is rather non-trivial, so we feel this is a hard problem.

Solution. For a contradiction, assume that

$$
\frac{1}{b_{1}}+\ldots+\frac{1}{b_{n}}>1
$$

So certainly $n \geqslant 2$. Note also that $A$ is infinite.
For each $i$, define $f_{i}(x)=b_{i} x+c_{i}$. Each function $f_{i}$ maps $A$ to itself. By condition (ii), if $f_{i}(a)=f_{j}\left(a^{\prime}\right)$ for $a, a^{\prime} \in A$ then $i=j$ and so $a=a^{\prime}$. By iterating this argument, it follows that if we have $a, a^{\prime} \in A$ and

$$
f_{i_{1}}\left(f_{i_{2}}\left(\ldots f_{i_{r}}(a) \ldots\right)\right)=f_{j_{1}}\left(f_{j_{2}}\left(\ldots f_{j_{r}}\left(a^{\prime}\right) \ldots\right)\right)
$$

then each $i_{k}=j_{k}\left(\right.$ and $\left.a=a^{\prime}\right)$.
The idea of our proof is to take a suitable $a \in A$ and generate a family of numbers of the form $f_{i_{1}}\left(f_{i_{2}}\left(\ldots f_{i_{r}}(a) \ldots\right)\right)$ by choosing $i_{1}, \ldots, i_{r}$ appropriately. These iterates are all distinct, and we will get an upper estimate on their size, but a lower bound on their number. Before embarking on the main argument, let us consider the case when $b_{1}=b_{2}=b_{3}=2$ as an illustration. If we choose $a \in A$ large enough, then for any $i_{k} \in\{1,2,3\}$ we have $f_{i_{1}}\left(f_{i_{2}}\left(\ldots f_{i_{r}}(a) \ldots\right)\right) \leqslant(2.01)^{r} a$. But there are $3^{r}$ such numbers, all distinct. Taking $r$ large enough we see this is impossible.
This argument is trickier to generalize when the $b_{i}$ are not all equal. To do so we will look at all sequences of the functions $f_{i}$ of length $N$ where each $f_{i}$ appears in a proportion dependent only on $i$. Take positive rational numbers $p_{1}, \ldots p_{n}$ (to be chosen later) with $p_{1}+\ldots+p_{n}=1$ and an integer $N$ which is a multiple of $N_{0}$, the least common multiple of the denominators of the $p_{i}$. Hence each $p_{i} N$ is an integer. Let $a \in A$ and consider the set $\Phi(N)$ of the numbers of the form $f_{i_{1}}\left(f_{i_{2}}\left(\ldots f_{i_{N}}(a) \ldots\right)\right)$ where, for each $i$, exactly $p_{i} N$ of the $i_{k}$ equal $i$. As these are all distinct, $\Phi(N)$ has

$$
|\Phi(N)|=\frac{N!}{\left(p_{1} N\right)!\ldots\left(p_{n} N\right)!}
$$

elements.
Choose $d_{i}>b_{i}$ with $1 / d_{1}+\ldots+1 / d_{n}>1$. There is a number $K$ such that $K \leqslant f_{i}(x) \leqslant d_{i} x$ for all $i$ whenever $x \geqslant K$. As $A$ is infinite, we may choose $a \in A$ with $a \geqslant K$. Then

$$
f_{i_{1}}\left(f_{i_{2}}\left(\ldots f_{i_{N}}(a) \ldots\right)\right) \leqslant d_{i_{1}} \ldots d_{i_{N}} a=d_{1}^{p_{1} N} \ldots d_{n}^{p_{N} N} a .
$$

To obtain a contradiction we need an estimate for $|\Phi(N)|$. We can use Stirling's formula, or similar, to get a precise asymptotic formula for $|\Phi(N)|$, but we can obtain a weaker but still adequate lower bound in a completely elementary fashion. Note that

$$
\begin{aligned}
& \frac{\left|\Phi\left(N+N_{0}\right)\right|}{|\Phi(N)|} \\
= & \frac{\left(N+N_{0}\right)\left(N+N_{0}-1\right) \ldots(N+1)}{\left[\left(p_{1} N+p_{1} N_{0}\right)\left(p_{1} N+p_{1} N_{0}-1\right) \ldots\left(p_{1} N+1\right)\right] \ldots\left[\left(p_{n} N+p_{n} N_{0}\right) \ldots\left(p_{n} N+1\right)\right]} \\
= & \frac{\left(1+N_{0} / N\right)\left(1+\left(N_{0}-1\right) / N\right) \ldots(1+1 / N)}{\left[\left(p_{1}+p_{1} N_{0} / N\right)\left(p_{1}+\left(p_{1} N_{0}-1\right) / N\right) \ldots\left(p_{1}+1 / N\right)\right] \ldots\left[\left(p_{n}+p_{n} N_{0} / N\right) \ldots\left(p_{n}+1 / N\right)\right]} .
\end{aligned}
$$

If we choose any $q>p_{1}^{p_{1}} \ldots p_{n}^{p_{n}}$ then $\left|\Phi\left(N+N_{0}\right)\right| /|\Phi(N)|>1 / q^{N_{0}}$ for large enough $N$. It follows that there is a constant $U$ with $|\Phi(N)|>U / q^{N}$ for all $N$ divisible by $N_{0}$. But, as the size of the elements of $\Phi(N)$ is at most $a\left(d_{1}^{p_{1}} \ldots d_{n}^{p_{n}}\right)^{N}$, we have $|\Phi(N)| \leqslant a\left(d_{1}^{p_{1}} \ldots d_{n}^{p_{n}}\right)^{N}$. It is now clear how to choose $p_{1}, \ldots, p_{n}$ and $q$. We take $p_{1}, \ldots, p_{n}$ proportional to $1 / d_{1}, \ldots, 1 / d_{n}$, and $q$ with $p_{1}^{p_{1}} \ldots p_{n}^{p_{n}}<q<1 /\left(d_{1}^{p_{1}} \ldots d_{n}^{p_{n}}\right)$. Then our bounds on $|\Phi(N)|$ are contradictory for large $N$.

## Further comment

2. Another approach is to consider the functions $N(x)=|\{a \in A: a \leqslant x\}|$ and $N_{i}(x)=\left|\left\{a \in b_{i} A+c_{i}: a \leqslant x\right\}\right|$. One gets

$$
N(x) \geqslant \sum_{i=1}^{n} N_{i}(x) \geqslant \sum_{i=1}^{n} N\left(\frac{x}{b_{i}}-u\right)
$$

for a suitable number $u$. By ingenious manipulation of this inequality, a contradiction can be obtained. However these manipulations seem less natural than the approach through iterating applications of the $f_{i}$.

C1. Let $n$ be a positive integer. Each point $(x, y)$ in the plane, where $x$ and $y$ are non-negative integers with $x+y<n$, is coloured red or blue, subject to the following condition: if a point $(x, y)$ is red, then so are all points $\left(x^{\prime}, y^{\prime}\right)$ with $x^{\prime} \leqslant x$ and $y^{\prime} \leqslant y$. Let $A$ be the number of ways to choose $n$ blue points with distinct $x$-coordinates, and let $B$ be the number of ways to choose $n$ blue points with distinct $y$-coordinates. Prove that $A=B$.

## Comment

1. This is an easy question, with an interesting variety of approaches. We give three different solutions: one is by induction on $n$, one is by induction on the number of red points, and one is a direct bijection.

Solution 1. Let the number of blue points with $x$-coordinate $i$ be $a_{i}$, and let the number of blue points with $y$-coordinate $i$ be $b_{i}$. Our task is to show that $a_{0} a_{1} \ldots a_{n-1}=b_{0} b_{1} \ldots b_{n-1}$, and to accomplish this we will show that $a_{0}, a_{1}, \ldots, a_{n-1}$ is a permutation of $b_{0}, b_{1}, \ldots, b_{n-1}$.
We prove this result by induction on $n$. The case $n=1$ is trivial, so we pass to the induction step: we may assume the result for all smaller values of $n$.


Consider first the case when every point $(x, y)$ with $x+y=n+1$ is blue. Ignoring these points, we have a configuration for $n-1$, with blue columns of sizes $a_{0}-1, a_{1}-1, \ldots, a_{n-2}-1$ and blue rows of sizes $b_{0}-1, b_{1}-1, \ldots, b_{n-2}-1$. It follows by the induction hypothesis that $a_{0}-1, a_{1}-1, \ldots, a_{n-2}-1$ is a permutation of $b_{0}-1, b_{1}-1, \ldots, b_{n-2}-1$, and since $a_{n-1}=b_{n-1}=1$ we are done.
Now suppose instead that some point $(k, n-1-k)$ is red. Then the entire rectangle of all points $(x, y)$ with $x \leqslant k$ and $y \leqslant n-1-k$ is red. Thus, considering just the points $(x, y)$ with $x<k$, the induction hypothesis tells us that $a_{0}, a_{1}, \ldots, a_{k-1}$ is a permutation of $b_{n-k}, b_{n-k+1}, \ldots, b_{n-1}$, and similarly we have that $a_{k+1}, a_{k+2}, \ldots, a_{n-1}$ is a permutation of $b_{0}, b_{1}, \ldots, b_{n-2-k}$. Since $a_{k}=b_{n-1-k}=0$, we are done.

Solution 2. As above, we wish to show that $a_{0}, a_{1}, \ldots, a_{n-1}$ is a permutation of $b_{0}, b_{1}, \ldots, b_{n-1}$. We prove this by induction on the number of red points: the result is trivial when there are no red points. Choose a red point $(x, y)$ with $x+y$ maximal. Then $a_{x}=b_{y}=n-1-x-y$. If we change this red point to blue, then we have a configuration with fewer red points, with all blue rows and columns unchanged except that the values of $a_{x}$ and $b_{y}$ decrease by 1 . So from the induction hypothesis we have that $a_{0}, a_{1}, \ldots, a_{n-1}$, with $a_{x}$ replaced by $a_{x}-1$, is a permutation of $b_{0}, b_{1}, \ldots, b_{n-1}$, with $b_{y}$ replaced by $b_{y}-1$. Since $a_{x}=b_{y}$, it follows that $a_{0}, a_{1}, \ldots, a_{n-1}$ is a permutation of $b_{0}, b_{1}, \ldots, b_{n-1}$, as required.

Solution 3. We give an explicit bijection between $a_{0}, a_{1}, \ldots, a_{n-1}$ and $b_{0}, b_{1}, \ldots, b_{n-1}$. If $a_{x}=0$ then also $b_{x}=0$, and we let $a_{x}$ correspond with $b_{x}$. If $a_{x}>0$, let $(x, y)$ be the bottom blue point in column $x$. Now, among the points $(x, y),(x-1, y+1),(x-2, y+2), \ldots$, there must be at least one that is the leftmost blue point of a row: let the first one be ( $x^{\prime}, y^{\prime}$ ). Then we let $a_{x}$ correspond with $b_{y^{\prime}}$.
This is clearly reversible: if $b_{y}>0$ then we let $(x, y)$ be the leftmost blue point in row $y$, choose the first point among $(x, y),(x+1, y-1),(x+2, y-2), \ldots$ that is the bottom blue point of a row, say $\left(x^{\prime}, y^{\prime}\right)$, and let $b_{y}$ correspond to $a_{x^{\prime}}$.


C2. For $n$ an odd positive integer, the unit squares of an $n \times n$ chessboard are coloured alternately black and white, with the four corners coloured black. A tromino is an $L$-shape formed by three connected unit squares. For which values of $n$ is it possible to cover all the black squares with non-overlapping trominos? When it is possible, what is the minimum number of trominos needed?

Solution. Write $n=2 m+1$. The key observation for the second part of the problem, which also helps in eliminating cleanly the case $n=5$ for the first part, is the following. Consider the black squares at an even height above the bottom row: there are $(m+1)^{2}$ of them, and no two are covered by any one tromino. So we always need at least $(m+1)^{2}$ trominos to cover. This disposes of the cases $n=3$ and $n=5$ (and $n=1$ ), as in each of these cases we have that $3(m+1)^{2}$ is greater than $n^{2}$, so that the black squares cannot be covered by trominos. It remains only to show that when $n \geqslant 7$ we can cover the black squares with exactly $(m+1)^{2}$ trominos. For $n=7$, the numbers make this just about possible, as $3(m+1)^{2}$ is 48 . There are several ways to achieve this. One simple way is to note that we can make a $2 \times 3$ rectangle from trominos, and two of these together make a $4 \times 3$ rectangle. If we lay four of these around the $7 \times 7$ board, we have covered every square except the central one. But now take a tromino that is adjacent to the central square: by the way the $4 \times 3$ rectangles have been built up, it may be moved to uncover a white square and cover the central black square instead.


The case $n=7$


$$
\begin{aligned}
& \text { From }(2 m-1) \times(2 m-1) \\
& \text { to }(2 m+1) \times(2 m+1)
\end{aligned}
$$

In general, having covered the black squares on a $(2 m-1) \times(2 m-1)$ board with $m^{2}$ trominos, let us form a $(2 m+1) \times(2 m+1)$ board by surrounding it with a $(2 m-1) \times 2$ rectangle and a $(2 m+1) \times 2$ rectangle. Now, we may break up the $(2 m-1) \times 2$ rectangle into $2 \times 2$ squares ( $m-2$ of them) and one $2 \times 3$ rectangle, so that its black squares may be covered by $(m-2)+2$ trominos, and similarly the black squares of the $(2 m+1) \times 2$ rectangle may be covered by $(m-1)+2$ trominos. This gives a total of $m^{2}+m+(m+1)$ trominos, as required.

## Comment

1. Covering problems are not uncommon, but this problem seems to be rather unusual. What is nice is that when one first thinks about the problem it seems very messy, but when one has made the key observation above it all becomes very clean and elegant!

C3. Let $n$ be a positive integer. A sequence of $n$ positive integers (not necessarily distinct) is called full if it satisfies the following condition: for each positive integer $k \geqslant 2$, if the number $k$ appears in the sequence then so does the number $k-1$, and moreover the first occurrence of $k-1$ comes before the last occurrence of $k$. For each $n$, how many full sequences are there?

Solution. We claim there are $n$ ! full sequences. To do this, we will construct a bijection with the set of permutations of $\{1,2, \ldots, n\}$.
Let $a_{1}, \ldots, a_{n}$ be a full sequence, and let $r=\max \left(a_{1}, \ldots, a_{n}\right)$. Then all the numbers from 1 to $r$ occur in $a_{1}, \ldots, a_{n}$. Let $S_{i}=\left\{k: a_{k}=i\right\}$ for $1 \leqslant i \leqslant r$. Then all the $S_{i}$ are nonempty, and they partition the set $\{1,2, \ldots, n\}$. The condition that the sequence is full means that $\min S_{k-1}<\max S_{k}$ for $2 \leqslant k \leqslant r$. Now we write down a permutation $b_{1}, \ldots, b_{n}$ of $\{1,2, \ldots, n\}$ by writing down the elements of $S_{1}$ in descending order, then the elements of $S_{2}$ in descending order and so on. This gives a map from full sequences to permutations of $\{1,2, \ldots, n\}$.
Note also that this map is reversible. Indeed, given a permutation $b_{1}, \ldots, b_{n}$ of $\{1,2, \ldots, n\}$ let $S_{1}=\left\{b_{1}, \ldots, b_{k_{1}}\right\}$ where $b_{1}>\ldots>b_{k_{1}}<b_{k_{1}+1}$, let $S_{2}=\left\{b_{k_{1}+1}, \ldots, b_{k_{2}}\right\}$ where $b_{k_{1}+1}>\ldots>b_{k_{2}}<b_{k_{2}+1}$ and so on. Then let $a_{j}=i$ whenever $i \in S_{j}$.
It follows that the full sequences are in bijection with the set of permutations of $\{1,2, \ldots, n\}$, as required.

## Comment

1. It is easy to guess, from some small examples, that the answer is $n$ !, but finding a bijection is not easy. An alternative proof goes by induction on $n$ : given a full sequence of length $n$, we form a sequence of length $n-1$ by removing from it the first occurrence of its highest number. It is easy to check that this sequence of length $n-1$ is full. One can then check that each full sequence of length $n-1$ arises in this way from exactly $n$ full sequences of length $n$.

C4. Let $T$ be the set of ordered triples $(x, y, z)$, where $x, y, z$ are integers with $0 \leqslant x, y, z \leqslant 9$. Players $A$ and $B$ play the following guessing game. Player $A$ chooses a triple $(x, y, z)$ in $T$, and Player $B$ has to discover $A$ 's triple in as few moves as possible. A move consists of the following: $B$ gives $A$ a triple ( $a, b, c$ ) in $T$, and $A$ replies by giving $B$ the number $|x+y-a-b|+|y+z-b-c|+|z+x-c-a|$. Find the minimum number of moves that $B$ needs to be sure of determining $A$ 's triple.

Solution. It is easy to see that two moves cannot be enough. Indeed, each answer is an even integer between 0 and 54 inclusive, so that there are 28 possibilities for each answer. Thus with two moves the number of possible outcomes is at most $28^{2}$, which is less than the required number of outcomes, namely 1000 .
We now set out to show that three moves are enough, by providing an explicit strategy. The first move should be $(0,0,0)$. The reply is $2(x+y+z)$, so that we now know the value of $s=x+y+z$. Clearly $0 \leqslant s \leqslant 27$, but to reduce the number of cases, in the algorithm below we may assume that $s \leqslant 13$. Indeed, if $s \geqslant 14$ then we perform the algorithm below, but always 'reflecting', i.e. asking $(9-a, 9-b, 9-c)$ instead of $(a, b, c)$ - we will recover the reflection of $(x, y, z)$ at the end.
Case 1: $s \leqslant 9$. This is the easy case. The second move should be $(9,0,0)$. We learn $y+z+(9-x-y)+(9-x-z)=18-2 x$, so we now know the value of $x$. And similarly asking $(0,9,0)$ tells us the value of $y$, so we are done (as $z=s-x-y$ ).
Case 2: $9<s \leqslant 13$. The second move should be ( $9, s-9$, 0 ). We learn $z+|9-x-z|+|9-x|$, which is say $2 k$, where $k=z$ if $x+z \geqslant 9$ and $k=9-x$ if $x+z<9$. Note that whichever value $k$ takes we do have $z \leqslant k \leqslant s$.
Case 2a: $s-k \leqslant 9$. The third move should be $(s-k, 0, k)$. We learn $y+|k-y-z|+|z-k|$. Since $k \leqslant y+z$ (if $k=z$ then this is obvious, while if $k=9-x$ then $k-y-z=9-s)$, this is just $y+(y+z-k)+(k-z)=2 y$. Thus we know $y$, and hence $x+z$. So we know whether $k$ is $z$ or $9-x$, and we are done.
Case 2b: $s-k>9$. The third move should be $(9, s-k-9, k)$. We learn
$|s-k-x-y|+|s-9-y-z|+\mid 9+k-x-z$, which is $(k-z)+(9-x)+(9+k-x-z)=18+2 k-2(x+z)$. So we know $x+z$, and so we know whether $k$ is $z$ or $9-x$. In either case, we know one of $x$ and $z$, and from $x+z$ we may deduce the other one, and we are done.

## Comments

1. Case 1 is the natural, simple case to deal with first. And then, in Case 2, the key idea is to use $s$ itself in the triple we ask. Similarly, in Case 2a the key idea is to somehow include $k$ in the triple we ask, while Case 2 b is just a more complicated version of Case 2a. One could view the second move in Case 2 as the obvious modification of the second move in Case 1 when $s$ is 'out of range', and similarly the third move in Case 2 b is the obvious modification of the third move in Case 2a.
2. As two moves only just fail, it is very natural to guess that three moves is the right answer. But finding an actual strategy seems complicated. This question requires clarity of thought, but no specialist knowledge at all.
3. There is nothing special about the fact that there are 10 possible values for each digit. Changing 10 to a larger number, such as 2002, would not change the solution.

C5. Let $r \geqslant 2$ be a fixed positive integer, and let $F$ be an infinite family of sets, each of size $r$, no two of which are disjoint. Prove that there exists a set of size $r-1$ that meets each set in $F$.

Solution. We will show the following: if $A$ is a set of size less than $r$ that is contained in infinitely many sets of $F$, then either $A$ meets all sets in $F$ (in which case we are done) or else there is an $x \notin A$ such that $A \cup\{x\}$ is itself contained in infinitely many sets of $F$. Since there certainly exists such a set $A$ (for example, the empty set), we may then iterate this result $r$ times and we will be done (as a set of size $r$ clearly cannot be contained in infinitely many sets of $F$ !). To prove the result, suppose that some set $R=\left\{x_{1}, x_{2}, \ldots, x_{r}\right\}$ in $F$ is disjoint from $A$. Of the infinitely many sets in $F$ that contain $A$, each must meet $R$, and so some $x_{i}$ is a member of infinitely many of them. So we may take $x=x_{i}$.

## Comments

1. Although the above proof is very short, it does seem to require a creative insight, namely the clever auxiliary result we prove.
2. An alternative proof is to note that, if the result is false, then for each set $R$ in $F$, and each point $x \in R$, the set $R-\{x\}$ is not suitable, so that there must be a set $S$ in $F$ with $R \cap S=\{x\}$. In other words, each point $x$ belonging to some set in $F$ is the intersection of two sets in $F$. This implies that there is no finite set $Y$ with the property that any two sets in $F$ must meet at some point belonging to $Y$. Although this is highly implausible, it does seem tricky to prove impossible. The simplest way is probably to prove a stronger result: that if $F$ and $G$ are families of $r$-sets, with each member of $F$ meeting each member of $G$, then there exists a finite set $Y$ such that each member of $F$ meets each member of $G$ at some point of $Y$. This goes easily by induction on $r$.
As in the proof above, this proof requires a key creative insight, namely the generalisation from one family to two. This does suggest that the problem may be quite hard.

C6. Let $n$ be an even positive integer. Show that there is a permutation $x_{1}, x_{2}, \ldots, x_{n}$ of $1,2, \ldots, n$ such that for every $1 \leqslant i \leqslant n$ the number $x_{i+1}$ is one of $2 x_{i}, 2 x_{i}-1,2 x_{i}-n$, $2 x_{i}-n-1$ (where we take $x_{n+1}=x_{1}$ ).

Solution. Write $n=2 m$. We shall define a directed graph $G$ with vertices labelled $1, \ldots, m$ and edges labelled $1, \ldots, 2 m$. The edges issuing from vertex $i$ are labelled $2 i-1$ and $2 i$, and those entering it are labelled $i$ and $i+m$.
The underlying graph of $G$ is connected: by induction on $j$ there is a path from 1 to $j$ since if $j>1$ then $j=2 k-1$ or $2 k$ with $1 \leqslant k<j$ and there is an edge from $k$ to $j$. Also, the indegree and outdegree of each vertex is the same (namely 2). The directed graph $G$ thus has an Euler circuit. Let $x_{i}$ be the label of the $i$-th edge in such a circuit. If edge $x_{i}$ enters and edge $x_{i+1}$ leaves vertex $j$ then $x_{i} \equiv j(\bmod m)$ and $x_{i+1}=2 j-1$ or $2 j$. Hence $2 x_{i} \equiv 2 j(\bmod 2 m=n)$ and so $x_{i+1} \equiv 2 x_{i}-1$ or $2 x_{i}(\bmod n)$ as is required.

## Comments

1. The problem requires one to prove the existence of a Hamilton cycle in a certain graph. There is no obvious way to do this. The above solution avoids this difficulty by constructing a directed graph on a smaller vertex set whose edges are labelled $1,2, \ldots, n$ in such a way that an Euler circuit (a closed walk that traverses each edge exactly once) corresponds to the desired permutation.
2. The above proof uses the simple fact that if all indegrees and outdegrees are equal (and the underlying graph is connected) then there is an Euler circuit. This is very easy to prove, for example by considering a closed walk of maximal length. Indeed, the proposer's solution does essentially the same thing with bare hands.
3. Although the above solution is short, it is definitely hard to think of. There is a slight similarity with the construction of de Bruijn sequences, but this is only apparent with hindsight once one has written down the proof.

C7. Among a group of 120 people, some pairs are friends. A weak quartet is a set of four people containing exactly one pair of friends. What is the maximum possible number of weak quartets?

Solution. We proceed in three steps. First, we will show that, for a maximum number of weak quartets, our graph (thinking of the friends as defining a graph on 120 vertices) breaks up as a disjoint union of complete graphs. Then we show that these complete graphs have sizes that are as equal as possible (ie. differ by at most 1 from each other). And lastly we will find which of these is the best.
For the first step, we would like to show that any two adjacent vertices have the same neighbours (apart from themselves, of course). For a graph $G$ on our 120 vertices, write $Q(G)$ for the number of weak quartets in $G$. For adjacent vertices $x, y$ of $G$, let $G^{\prime}$ be the graph formed from $G$ by 'copying' $y$ to $x$ : in other words, for each $z \neq x, y$, we add the edge $x z$ if $y z$ is an edge and we remove the edge $x z$ if $y z$ is not an edge. Similarly, let $G^{\prime \prime}$ be the graph formed from $G$ by copying $x$ to $y$.
We claim that $Q(G) \leqslant \frac{1}{2}\left(Q\left(G^{\prime}\right)+Q\left(G^{\prime \prime}\right)\right.$. Indeed, let us compare the weak quartets in $G$ with those in $G^{\prime}$ and $G^{\prime \prime}$. The number of weak quartets containing neither $x$ nor $y$ is the same in $G, G^{\prime}$ and $G^{\prime \prime}$, while the number containing both $x$ and $y$ is at least as great in $G^{\prime}$ and $G^{\prime \prime}$ as it is in $G$. The number containing $y$ but not $x$ in $G^{\prime}$ is at least twice what it is in $G$, while the number containing $x$ but not $y$ in $G^{\prime \prime}$ is at least twice what it is in $G$. This establishes our claim.
It follows that, for an extremal $G$, we must have $Q(G)=Q\left(G^{\prime}\right)=Q\left(G^{\prime \prime}\right)$. So we may repeat this copying operation pair by pair, to obtain a graph in which any two adjacent vertices have the same common neighbours. Indeed, if $x$ and $y$ are two adjacent vertices then we copy $y$ to $x$; if there another vertex $z$ adjacent to $x$ (and so also to $y$ ) then we copy $z$ to $x$ and then to $y$, and so on. This completes the first step.
Our aim now is to show that the sizes of the complete graphs in $G$ may be taken to be as equal as possible. There are various ways to do this: one way is as follows. Let the complete graphs in $G$ have sizes $a_{1}, a_{2}, \ldots, a_{n}$, where, just for convenience in what is to follow, we allow $a_{i}=0$. Then we have

$$
Q(G)=\sum_{i=1}^{n}\binom{a_{i}}{2} \sum_{j<k, j, k \neq i} a_{j} a_{k} .
$$

(Here as usual we just take $\binom{a}{2}$ to mean $a(a-1) / 2$, to cover the case when $a<2$.) Now, let us consider two of the $a_{i}$, say $a$ and $b$, and let us see how $Q(G)$ varies as we change the values of $a$ and $b$ (keeping the other values, and the sum $a+b$, fixed). We have

$$
Q(G)=A\left(\binom{a}{2}+\binom{b}{2}\right)+B(a+b)+C\left(\binom{a}{2} b+\binom{b}{2} a\right),
$$

where $A, B, C$ do not depend on $a$ or $b$. If we swap $a$ and $b$, we get the same expression, which tells us that the expression is a quadratic, symmetric about $s / 2$, where $s$ is the fixed sum $a+b$. (The expression may appear to be cubic, but it is easy to see that there is no cubic term, either by direct calculation or because no cubic can be symmetric!)
This tells us that the maximum when $0 \leqslant a \leqslant s$ occurs either at $a=s / 2$ or at $a=0$ (and $a=s$ ), and that the maximum value for integer $a$ occurs when $a=0$ or when $a=b$ or when $a=b \pm 1$. Repeating for each pair of the $a_{i}$, we have completed the second step.

Our final step is some calculation. Writing $n$ for the number of (non-empty) complete graphs, we see that

$$
Q(G)=n\binom{120 / n}{2}\binom{n-1}{2}(120 / n)^{2}
$$

whenever $n$ divides 120 . It is easy to check that, for $n \leqslant 6$, the maximum occurs at $n=5$, with value $15.23 .24^{3}$. Moreover, because of the fact that the maximum over all real $a_{i}$ in the previous paragraph occurred when all the non-zero $a_{i}$ were equal, we also know that the maximum possible value of $Q(G)$ is at most the maximum value of the expression

$$
n\binom{120 / n}{2}\binom{n-1}{2}(120 / n)^{2}
$$

as $n$ varies from 3 to 120 . But this function is at most $120^{4}(n-1)(n-2) / 4 n^{3}$, which is a decreasing function of $n$ for $n \geqslant 6$ and is at most $15.23 .24^{3}$ for $n=7$. This completes the third step: the maximum value is $15.23 .24^{3}$.

## Comments

1. The rough strategy outlined at the start of the proof is not too hard to think of. However, the actual extremal configuration (with 5 complete graphs of size 24) is far from obvious. In addition, fitting in the detail to implement the outline strategy presents a number of challenges. The hardest of these is the first step, to show that we have a disjoint union of complete graphs.
2. The number 120 has been chosen to make the numerical calculations as simple as possible.
