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On some stability concepts for real functions

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Abstract. In this paper we investigate some asymptotic properties for real functions defined on \mathbb{R}_+ as exponential stability, uniform exponential stability, polynomial stability and uniform polynomial stability. Our main objectives are to give characterizations for these concepts and to establish connections between them.

Keywords: Exponential stability, polynomial stability. **MSC** : 34D05

1. Exponential stability

The exponential stability property plays a central role in the theory of asymptotic behaviors for differential equations. In this section we consider two concepts of exponential stability for the particular case of real functions defined on \mathbb{R}_+ .

Let Δ_0 be the set defined by

$$\Delta_0 = \{ (y, x) \in \mathbb{R}^2_+ \text{ with } y \ge x \}.$$

Definition 1. A real function $F : \mathbb{R}_+ \to \mathbb{R}$ is called

(i) uniformly exponentially stable (and denote u.e.s), if there are $N \ge 1$ and $\alpha > 0$ such that

$$|F(y)| \leq Ne^{-\alpha(y-x)}|F(x)|, \text{ for all } (y,x) \in \Delta_0;$$

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(ii) (nonuniformly) exponentially stable (and denote e.s), if there exist $\alpha > 0$ and a nondecreasing function $N : \mathbb{R}_+ \to [1, \infty)$ such that

$$|F(y)| \le N(x)e^{-\alpha(y-x)}|F(x)|, \text{ for all } (y,x) \in \Delta_0.$$

Remark 1. It is obvious that $u.e.s \Rightarrow e.s.$ The converse implication is not true, as shown in

Example 1. Let $F : \mathbb{R}_+ \to \mathbb{R}$ be the function defined by

$$F(x) = e^{x \cos x - 3x}$$

We observe that for all $(y, x) \in \Delta_0$, we have

$$|F(y)| = F(y) = e^{y \cos y - 3y - x \cos x + 3x} |F(x)| \le e^{-2y + 4x} |F(x)| \le e^{2x} e^{-2(y-x)} |F(x)|$$

and hence F is e.s.

If we suppose that if F is u.e.s, then there exist $N \geq 1$ and $\alpha > 0$ such that

$$e^{y\cos y-3y} \leq Ne^{-\alpha(y-x)}e^{x\cos x-3x}$$
, for all $(y,x) \in \Delta_0$

In particular, for $y = 2n\pi$ and $x = 2n\pi - \pi/2$, we obtain

$$e^{2n\pi} < N e^{(3-\alpha)\pi/2},$$

which for $n \to \infty$, leads to a contradiction.

A necessary and sufficient condition for uniform exponential stability is given by

Proposition 1. A function $F : \mathbb{R}_+ \to \mathbb{R}$ is uniformly exponentially stable if and only if there exists a decreasing function $f : \mathbb{R}_+ \to \mathbb{R}^*_+ = (0, \infty)$ with $\lim_{x\to\infty} f(x) = 0$ and

$$|F(y)| \leq f(y-x)|F(x)|, \text{ for all } (y,x) \in \Delta_0.$$

Proof. Necessity. It is an immediate verification.

Sufficiency. Let $\delta > 1$ with $f(\delta) < 1$. Then for all $(y, x) \in \Delta_0$ there exist $n \in \mathbb{N}$ and $r \in [0, \delta)$ such that $y = x + n\delta + r$. Then

$$\begin{aligned} |F(y)| &\leq f(r)|F(x+n\delta)| \leq f(0)|F(x+n\delta)| \leq \\ &\leq f(0)f(\delta)|F(x+(n-1)\delta)| \leq \dots \leq f(0)f^n(\delta)|F(x)| = \\ &= f(0)e^{n\ln f(\delta)}|F(x)| = f(0)e^{\alpha r}e^{-\alpha(y-x)}|F(x)| \leq \\ &\leq f(0)e^{\alpha\delta}e^{-\alpha(y-x)}|F(x)| \leq Ne^{-\alpha(y-x)}|F(x)|, \end{aligned}$$

where $\alpha = -\frac{\ln \delta}{f(\delta)}$ and $N = 1 + f(0)e^{\alpha\delta}$.

Definition 2. A function $F : \mathbb{R}_+ \to \mathbb{R}$ with the property that there are $M \ge 1$ and $\omega > 0$ such that

$$|F(y)| \le M e^{\omega(y-x)} |F(x)|, \text{ for all } (y,x) \in \Delta_0,$$

is called with exponential growth (and we denote e.g).

Remark 2. It is obvious that $u.e.s \Rightarrow e.g$ and the converse implication is not true.

Proposition 2. A function $F : \mathbb{R}_+ \to \mathbb{R}$ is with exponential growth if and only if there exists a nondecreasing function $\varphi : \mathbb{R}_+ \to \mathbb{R}^*_+$ with $\lim_{x\to\infty} \varphi(x) = \infty$ and

$$|F(y)| \le \varphi(y-x)|F(x)|, \text{ for all } (y,x) \in \Delta_0.$$

Proof. It is similar with the proof of Proposition 1.

Another characterization of the u.e.s property is given by

Proposition 3. Let $F : \mathbb{R}_+ \to \mathbb{R}$ be an integrable function on each compact interval $[a,b] \subset \mathbb{R}_+$ (i.e., F is locally integrable on \mathbb{R}_+) with exponential growth. Then F is uniformly exponentially stable if and only if there exists $D \ge 1$ with

$$\int_{x}^{\infty} |F(y)| \mathrm{d}y \le D|F(x)|, \text{ for all } x \in \mathbb{R}_{+}.$$

Proof. Necessity. It is an immediate verification.

Sufficiency. Because F is with e.g, it follows that there is a nondecreasing function $\varphi : \mathbb{R}_+ \to [1, \infty)$ with

$$|F(y)| \le \varphi(y-x)|F(x)|$$
, for all $(y,x) \in \Delta_0$.

Then for all $(y, x) \in \Delta_0$ with $y \ge x + 1$ we have

$$|F(y)| = \int_{y-1}^{y} |F(y)| dz \le \int_{y-1}^{y} \varphi(y-z)|F(z)| dz \le$$
$$\le \varphi(1) \int_{x}^{\infty} |F(z)| dz \le D\varphi(1)|F(x)|.$$

If $(y, x) \in \Delta_0$ with $y \in [x, x+1)$, then

$$|F(y)| \le \varphi(y-x)|F(x)| \le D\varphi(1)|F(x)|$$

and hence

$$|F(y)| \leq D\varphi(1)|F(x)|$$
, for all $(y, x) \in \Delta_0$.

Thus we obtain

$$(y-x)|F(y)| = \int_{x}^{y} |F(y)| dz \le D\varphi(1) \int_{x}^{y} |F(z)| dz \le D^{2}\varphi(1)|F(x)|,$$

for all $(y, x) \in \Delta_0$ and

$$|F(y)| \le f(y-x)|F(x)|$$
, for all $(y,x) \in \Delta_0$,

where

$$f(z) = \frac{D^2\varphi(1)}{z+1}$$

Using the preceding proposition we conclude that F is u.e.s. \Box A generalization of the preceding result for the nonuniform case is given

by

Proposition 4. Let $F : \mathbb{R}_+ \to \mathbb{R}$ be a locally integrable function on \mathbb{R}_+ with exponential growth. Then F is exponentially stable if and only if there are $\beta > 0$ and $D : \mathbb{R}_+ \to [1, \infty)$ such that

$$\int_{x}^{\infty} e^{\beta y} |F(y)| \mathrm{d}y \le D(x) e^{\beta x} |F(x)|, \text{ for all } x \ge 0.$$

Proof. Necessity. It is a simple verification for $\beta \in (0, \alpha)$, where α is given by Definition 1 (ii).

Sufficiency. If $(y, x) \in \Delta_0$ with $y \ge x + 1$ then

$$\begin{split} e^{\beta y}|F(y)| &= \int_{y-1}^{y} e^{\beta y}|F(y)| \mathrm{d}z \leq \int_{y-1}^{y} \varphi(y-z) e^{\beta y}|F(z)| \leq \\ &\leq \varphi(1) \int_{y-1}^{y} e^{\beta(y-z)} e^{\beta z}|F(z)| \mathrm{d}z \leq \varphi(1) e^{\beta} \int_{x}^{\infty} e^{\beta z}|F(z)| \mathrm{d}z \leq \\ &\leq D(x) e^{\beta} \varphi(1) e^{\beta x}|F(x)|. \end{split}$$

If $(y, x) \in \Delta_0$ with $y \in [x, x+1)$ then

$$e^{\beta y}|F(y)| \le e^{\beta(y-x)}e^{\beta x}\varphi(y-x)|F(x)| \le \\ \le e^{\beta}\varphi(1)e^{\beta x}|F(x)| \le D(x)e^{\beta}\varphi(1)e^{\beta x}|F(x)|.$$

Finally, we obtain

$$|F(y)| \le N(x)e^{-\beta(y-x)}|F(x)|, \text{ for all } (y,x) \in \Delta_0,$$

where

$$N(x) = \varphi(1)e^{\beta}D(x).$$

In the particular case of u.e.s property we obtain:

Corollary 1. Let $F : \mathbb{R}_+ \to \mathbb{R}$ be a locally integrable function on \mathbb{R}_+ with exponential growth. Then F is uniformly exponentially stable if and only if there exist $D \ge 1$ and $\beta > 0$ such that

$$\int_{x}^{\infty} e^{\beta y} |F(y)| \mathrm{d}y \le D e^{\beta x} |F(x)|, \text{ for all } x \ge 0.$$

2. Polynomial stability

In this section we consider two concepts of polynomial stability for real functions. Our approach is based on the extension of techniques of exponential stability to the case of polynomial stability. Two illustrating examples clarify the relations between the stability concepts considered in this paper.

Let Δ_1 be the set defined by

$$\Delta_1 = \{(y, x) \in \mathbb{R}^2 \text{ with } y \ge x \ge 1\}.$$

Definition 3. A real function $F : \mathbb{R}_+ \to \mathbb{R}$ is called

 (i) uniformly polynomially stable (and denote u.p.s) if there exist N ≥ 1 and α > 0 such that

$$y^{\alpha}|F(y)| \leq Nx^{\alpha}|F(x)|, \text{ for all } (y,x) \in \Delta_1;$$

(ii) (nonuniformly) polynomially stable (and denote p.s) if there are a nondecreasing function $N : \mathbb{R}_+ \to [1, \infty)$ and $\alpha > 0$ such that

 $y^{\alpha}|F(y)| \leq N(x)x^{\alpha}|F(x)|, \text{ for all } (y,x) \in \Delta_1.$

Remark 3. If the function $F : \mathbb{R}_+ \to \mathbb{R}$ is u.e.s then it is u.p.s. Indeed, if *F* is u.e.s then using the monotony of the function

$$f: [1,\infty) \to [e,\infty), \ f(t) = \frac{e^t}{t}$$

we obtain

$$y^{\alpha}|F(y)| \le x^{\alpha}e^{\alpha(y-x)}|F(x)| \le Nx^{\alpha}|F(x)|$$

for all $(y, x) \in \Delta_1$ and hence F is u.p.s. The converse is not true, phenomenon illustrated by

Example 2. The function $F : \mathbb{R}_+ \to \mathbb{R}$, $F(x) = \frac{1}{x^3 + 1}$, satisfies the inequality

$$|y^{3}|F(y)| = \frac{y^{3}}{y^{3}+1} \le \frac{2x^{3}}{x^{3}+1} = 2x^{3}|F(x)|,$$

for all $(y,x) \in \Delta_1$. This shows that F is u.p.s. We show that F is not e.s and hence it is not u.e.s. Indeed, if we suppose that F is e.s then there are

 $\alpha > 0$ and $N : \mathbb{R}_+ \to [1, \infty)$ such that

$$\frac{e^{\alpha y}}{y^3+1} \le N(x)\frac{e^{\alpha x}}{x^3+1}, \text{ for all } (y,x) \in \Delta_0.$$

For x = 0 and $y \to \infty$ we obtain a contradiction which proves that F is not e.s.

Remark 4. Similarly, as in Remark 3 we can prove that if $F : \mathbb{R}_+ \to \mathbb{R}$ is *e.s then it is p.s.*

The function considered in Example 2 shows that the converse is not true.

Remark 5. It is obvious that if the function $F : \mathbb{R}_+ \to \mathbb{R}$ is u.p.s then it is p.s. The function F considered in Example 1 shows that the converse implication is not true. Indeed, F is e.s. and by Remark 4 it is p.s.

If we suppose that F is u.p.s then there are $\alpha > 0$ and $N \ge 1$ such that

$$y^{\alpha}e^{y\cos y-3y} \leq Nx^{\alpha}e^{x\cos x-3x}, \text{ for all } (y,x) \in \Delta_1$$

Then for $x = 2n\pi - \pi/2$, $y = 2n\pi$ and $n \to \infty$ we obtain a contradiction.

Remark 6. The preceding considerations prove that between the asymptotic behaviors defined in this paper we have the following implications:

$$\begin{array}{cccc} e.g & \Leftarrow & u.e.s \Rightarrow & e.s \\ \Downarrow & & \Downarrow & & \Downarrow \\ p.g & \Leftarrow & u.p.s \Rightarrow & p.s \end{array}$$

Definition 4. A function $F : \mathbb{R}_+ \to \mathbb{R}$ with the property that there are $M \ge 1$ and p > 0 such that

 $x^p|F(y)| \leq My^p|F(x)|, \text{ for all } (y,x) \in \Delta_1$

is called with polynomial growth (and denote p.g).

A necessary and sufficient condition for polynomial stability is given by

Proposition 5. Let $F : \mathbb{R}_+ \to \mathbb{R}$ be a locally integrable function on \mathbb{R}_+ with polynomial growth. Then F is polynomially stable with $\alpha > 1$ if and only if there exists $\beta > 0$ and $D : [1, \infty) \to [1, \infty)$ such that

$$\int_{x}^{\infty} y^{\beta} |F(y)| \mathrm{d}y \le D(x) x^{\beta+1} |F(x)|, \text{ for all } x \ge 1.$$

Proof. Necessity. If F is p.s with $\alpha > 1$ then for $\beta \in (0, \alpha - 1)$ we have

$$\int_{x}^{\infty} y^{\beta} |F(y)| \mathrm{d}y \le N(x) |F(x)| x^{\alpha} \int_{x}^{\infty} y^{\beta-\alpha} \mathrm{d}y \le$$
$$\le \frac{N(x)}{\alpha - 1 - \beta} x^{\beta+1} |F(x)| \le D(x) x^{\beta+1} |F(x)|, \text{ for all } x \ge 1$$

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where

$$D(x) = 1 + \frac{N(x)}{\alpha - 1 - \beta}.$$

Sufficiency. If $(y, x) \in \Delta_1$ and $y \in [x, 2x)$ then

$$y^{\beta+1}|F(y)| \le My^{\beta+1} \left(\frac{y}{x}\right)^p |F(x)| \le Mx^{\beta+1} \left(\frac{y}{x}\right)^{p+\beta+1} |F(x)| \le M2^{p+\beta+1} x^{\beta+1} |F(x)|,$$

where M and p are given by Definition 4.

If $(y, x) \in \Delta_1$ with $y \ge 2x$ and

$$C = \int_{1}^{2} \frac{\mathrm{d}z}{z^{p+\beta+2}} = \frac{1}{y^{p+\beta+1}} \int_{\frac{y}{2}}^{y} t^{p+\beta} \mathrm{d}t,$$

then

$$Cy^{\beta+1}|F(y)| = \frac{1}{y^p} \int_{\frac{y}{2}}^{y} t^{p+\beta}|F(y)| \mathrm{d}t \le$$
$$\le M \int_{x}^{\infty} t^{\beta}|F(t)| \mathrm{d}t \le MD(x)x^{\beta+1}|F(x)|$$

Finally, we obtain

$$y^{\alpha}|F(y)| \leq N(x)x^{\alpha}|F(x)|, \text{ for all } (y,x) \in \Delta_1,$$

where $\alpha = \beta + 1 > 1$ and

$$N(x) = M(2^{p+\beta+1} + D(x)) \ge 1,$$

which shows that F is p.s. with $\alpha > 1$.

From the proof of the preceding result we obtain

Corollary 2. Let $F : \mathbb{R}_+ \to \mathbb{R}$ be a locally integrable function on \mathbb{R}_+ with polynomial growth. Then F is uniformly polynomially stable with $\alpha > 1$ if and only if there are $D \ge 1$ and $\beta > 0$ such that

$$\int_{x}^{\infty} y^{\beta} |F(y)| \mathrm{d}y \le Dx^{\beta+1} |F(x)|, \text{ for all } x \ge 1.$$

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Some inequalities about certain arithmetic functions which use the e-divisors and the e-unitary divisors

NICUŞOR MINCULETE¹⁾

Abstract. The purpose of this paper is to present several inequalities about the arithmetic functions $\sigma^{(e)}$, $\tau^{(e)}$, $\sigma^{(e)*}$, $\tau^{(e)*}$ and other well-known arithmetic functions. Among these, we have the following: $\frac{\tau(n)}{\tau^*(n)} \geq \frac{\tau^{(e)}(n)}{\tau^{(e)*}(n)}, \frac{\sigma(n)}{\sigma^*(n)} \geq \frac{\sigma^{(e)}(n)}{\sigma^{(e)*}(n)}, \tau(n) + 1 \geq \tau^{(e)}(n) + \tau^*(n)$ and $\sigma(n) + n \geq \sigma^{(e)}(n) + \sigma^*(n)$, for any $n \geq 1$, where $\tau(n)$ is the number of natural divisors of n, $\tau^*(n)$ is the number of natural divisors of n, $\tau^*(n)$ is the number of unitary divisors of n, $\sigma^*(n)$ is the sum of the unitary divisors of n and γ is the "core" of n. **Keywords:** arithmetic function, exponential divisor, exponential unitary divisor.

MSC : 11A25

1. Introduction

First we have to mention that the notion of "exponential divisor" was introduced by M. V. Subbarao in [9], in the following way: if n > 1 is an integer of canonical form $n = p_1^{a_1} p_2^{a_2} \cdots p_r^{a_r}$, then the integer $d = \prod_{i=1}^r p_i^{b_i}$ is called an *exponential divisor* (or *e-divisors*) of $n = \prod_{i=1}^r p_i^{a_i} > 1$, if $b_i | a_i$ for every $i = \overline{1, r}$. We note $d|_{(e)}n$. Let $\sigma^{(e)}(n)$ denote the sum of the exponential divisors of n and $\tau^{(e)}(n)$ denote the number of exponential divisors of n.

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For example, if $n = 2^4 3^2$, then the exponential divisors of n are the following:

 $2 \cdot 3, 2 \cdot 3^2, 2^2 \cdot 3, 2^2 \cdot 3^2, 2^4 \cdot 3$ and $2^4 \cdot 3^2$.

For various properties of the arithmetic functions which use the edivisors see the monograph of J. Sándor and B. Crstici [5].

J. Fabrykowski and M. V. Subbarao in [1] study the maximal order and the average order of the multiplicative function $\sigma^{(e)}(n)$. E.G. Straus and M. V. Subbarao in [8] obtained several results concerning e-perfect numbers (nis an e-perfect number if $\sigma^{(e)}(n) = 2n$). They conjecture that there is only a finite number of e-perfect numbers not divisible by any given prime p.

In [4], J. Sándor showed that, if n is a perfect square, then

$$2^{\omega(n)} \le \tau^{(e)}(n) \le 2^{\Omega(n)},\tag{1}$$

where $\omega(n)$ and $\Omega(n)$ denote the number of distinct prime factors of n, and the total number of prime factors of n, respectively. It is easy to see that, for $n = p_1^{a_1} p_2^{a_2} \cdots p_r^{a_r} > 1$, we have $\omega(n) = r$ and $\Omega(n) = a_1 + a_2 + \ldots + a_r$.

In [6], J. Sándor and L. Tóth proved the inequality

$$\frac{n^k+1}{2} \ge \frac{\sigma_k^*(n)}{\tau^*(n)} \ge \sqrt{n^k},\tag{2}$$

for all $n \ge 1$ and $k \ge 0$, where $\tau^*(n)$ is the number of the unitary divisors of $n, \sigma_k^*(n)$ is the sum of kth powers of the unitary divisors of n.

In [11] *L. Tóth* and *N. Minculete* presented the notion of "exponential unitary divisors" or "e-unitary divisors". The integer $d = \prod_{i=1}^{r} p_i^{b_i}$ is called an *e-unitary divisor* of $n = \prod_{i=1}^{r} p_i^{a_i} > 1$ if b_i is a unitary divisor of a_i , so

 $\left(b_i, \frac{a_i}{b_i}\right) = 1$, for every $i = \overline{1, r}$. Let $\sigma^{(e)*}(n)$ denote the sum of the e-unitary divisors of n, and $\tau^{(e)*}(n)$ denote the number of the e-unitary divisors of n.

For example, if $n = 2^4 3^2$, then the exponential unitary divisors of n are the following:

$$2 \cdot 3, \ 2 \cdot 3^2, \ 2^4 \cdot 3 \text{ and } 2^4 \cdot 3^2.$$

By convention, 1 is an exponential divisor of itself, so that

$$\sigma^{(e)*}(1) = \tau^{(e)*}(1) = 1.$$

We notice that 1 is not an e-unitary divisor of n > 1, the smallest e-unitary divisor of $n = p_1^{a_1} p_2^{a_2} \cdots p_r^{a_r} > 1$ is $p_1 p_2 \cdots p_r = \gamma(n)$ is called the "core" of n.

In [3], it is show that

$$\sigma^{(e)}(n) \le \psi(n) \le \sigma(n),\tag{3}$$

where $\psi(n)$ is the function of Dedekind, and

$$\tau(n) \le \frac{\sigma^{(e)}(n)}{\tau^{(e)}(n)},\tag{4}$$

for all integers $n \ge 1$.

Other properties of the sum of the exponential divisors of n and of the number of the exponential divisors of n can be found in the papers [2, 6 and 10].

2. Inequalities for the functions $\tau^{\rm e}, \sigma^{\rm e}, \tau^{\rm e*}$ and $\sigma^{\rm e*}$

Lemma 1. For every $n \ge 1$, there is the following inequality

$$n\tau^*(n) \ge \sigma(n),$$
 (5)

and for all $n = p_1^{a_1} p_2^{a_2} \cdots p_r^{a_r} > 1$, with $a_i \ge 2$, (\forall) $i = \overline{1, r}$, there is the following inequality

$$n\tau^{\mathbf{e}}(n) \ge \sigma(n). \tag{6}$$

Proof. For $a \ge 1$ we show that

$$2p^a \ge p^a + p^{a-1} + \dots + p + 1.$$

This inequality is rewritten as

$$p^a \ge p^{a-1} + \dots + p + 1,$$

which by multiplication with p-1 becomes

$$p^{a+1} + 1 \ge 2p^a,$$

which is true, because $p^{a+1} \ge 2p^a$, for every prime number p and for all $a \ge 1$.

Therefore, $p^a \tau^*(p^a) = 2p^a \ge p^a + p^{a-1} + \ldots + p + 1 = \sigma(p^a)$, so $p^a \tau^*(p^a) \ge \sigma(p^a)$. Because the arithmetic functions τ^* and σ are multiplicative, we deduce the inequality

$$n\tau^*(n) \ge \sigma(n)$$

Since $a \ge 2$, it follows that $\tau(a) \ge 2$, so

$$p^{a}\tau(a) \ge 2p^{a} \ge p^{a} + \dots + p + 1 = \sigma(p^{a}),$$

which is equivalent to $p^a \tau^{e}(p^a) = p^a \tau(a) \ge \sigma(p^a)$. As the arithmetic functions τ^{e} and σ are multiplicative, we get the inequality of the statement. \Box

Remark 1. It may be noted that the lemma readily implies the inequality

$$n^2 \tau^{\mathbf{e}}(n^2) \ge \sigma(n^2),\tag{7}$$

for all $n \geq 1$.

Theorem 2. For every $n \ge 1$, there is the inequality

$$\frac{\tau(n)}{\tau^*(n)} \ge \frac{\tau^{\mathrm{e}}(n)}{\tau^{\mathrm{e}*}(n)}.$$
(8)

Proof. For n = 1, we have $\frac{\tau(1)}{\tau^*(1)} = 1 = \frac{\tau^{e}(1)}{\tau^{e*}(1)}$.

Let's consider $n = p_1^{a_1} p_2^{a_2} \cdots p_r^{a_r}$, with $a_i \ge 1$, $(\forall) i = \overline{1, r}$. According to lemma 1, we deduce the relation

$$a_i \tau^*(a_i) \ge \sigma(a_i), \ (\forall) \ i = \overline{1, r}.$$
 (9)

From the inequality of S. Sivaramakrishnan and C. S. Venkataraman [7], $\sigma(n) \ge \sqrt{n\tau(n)}, (\forall) \ n \ge 1$, we get $\sigma(a_i) \ge \sqrt{a_i\tau(a_i)}, (\forall) \ i = \overline{1, r}$. This last inequality combined with inequality (9) implies the inequality

$$\sqrt{a_i} \geq \frac{\tau(a_i)}{\tau^*(a_i)}, \ (\forall) \ i = \overline{1, r}, \tag{10}$$

which means that

$$\sqrt{\prod_{i=1}^{r} a_i} \ge \frac{\tau^{\mathrm{e}}(n)}{\tau^{\mathrm{e}*}(n)}, \text{ for all } n = p_1^{a_1} p_2^{a_2} \cdots p_r^{a_r} > 1, \text{ because}$$
(11)

 $\begin{aligned} \tau^{\mathrm{e}}(n) &= \tau(a_1) \cdot \ldots \cdot \tau(a_r) \text{ and } \tau^{\mathrm{e*}}(n) = \tau^{*}(a_1) \cdot \ldots \cdot \tau^{*}(a_r).\\ \text{But } a_i + 1 \geq 2\sqrt{a_i}, \ (\forall) \ i = \underbrace{1, r}_{i}, \text{ and by taking the product, we obtain} \end{aligned}$

the inequality $\prod_{i=1}^{r} (a_i + 1) \ge 2^r \sqrt{\prod_{i=1}^{r} a_i}$, which is equivalent to the relation $\tau(n) \ge \tau^*(n) \sqrt{\prod_{i=1}^{r} a_i}$, so

$$) \ge \tau^{*}(n) \sqrt{\prod_{i=1}^{r} a_{i}}, \text{ so}$$
$$\frac{\tau(n)}{\tau^{*}(n)} \ge \sqrt{\prod_{i=1}^{r} a_{i}} \text{ for all } n = p_{1}^{a_{1}} p_{2}^{a_{2}} \dots p_{r}^{a_{r}} > 1.$$
(12)

Combining relations (11) and (12), we deduce inequality (8). Finally, the proof is completed. $\hfill \Box$

Theorem 3. For every $n \ge 1$, there is the equality

$$\frac{\sigma(n)}{\sigma^*(n)} \ge \frac{\sigma^{\rm e}(n)}{\sigma^{\rm e*}(n)}.$$
(13)

Proof. We distinguish the following cases: Case I. For n = 1, we have $\frac{\sigma(1)}{\sigma^*(1)} = 1 = \frac{\sigma^{\mathrm{e}}(1)}{\sigma^{\mathrm{e}*}(1)}$. Case II. If n is squarefree, then $\sigma(n) = \sigma^*(n)$, and $\sigma^{\mathrm{e}}(n) = n = \sigma^{\mathrm{e}*}(n)$. Therefore, we obtain the relation of statement. Case III. Let's consider $n = p_1^{a_1} p_2^{a_2} \cdots p_r^{a_r} > 1$, with $a_i \ge 2$, $(\forall) i = \overline{1, r}$, then

$$\frac{\sigma(p^a)}{\sigma^*(p^a)} = \frac{1+p+p^2+\ldots+p^a}{1+p^a} = 1 + \frac{p+p^2+\ldots+p^{a-1}}{1+p^a}$$
(14)

and, because the exponential unitary divisors $p^{d_1}, ..., p^{d_q}$ of p^a are among the exponential divisors p^{d_1}, \dots, p^{d_s} of p^a , we have the inequality

$$\frac{\sigma^{\mathbf{e}}(p^{a})}{\sigma^{\mathbf{e}*}(p^{a})} = \frac{p^{d_{1}} + p^{d_{2}} + \dots + p^{d_{s}}}{p^{d_{1}} + \dots + p^{d_{q}}} = 1 + \frac{p^{d_{2}} + \dots + p^{d_{s-1}}}{p + \dots + p^{a}} \le 1 + \frac{p^{2} + p^{3} + \dots + p^{a-1}}{p + p^{a}} \le 1 + \frac{p + \dots + p^{a-1}}{1 + p^{a}}.$$

Therefore, using the inequality (14), we get the relation

$$\frac{\sigma^{\mathbf{e}}(p^a)}{\sigma^{\mathbf{e}*}(p^a)} \le \frac{\sigma(p^a)}{\sigma^*(p^a)}$$

Because the arithmetic functions σ^{e} , σ^{e*} , σ and σ^{*} are multiplicative, we deduce the inequality

$$\frac{\sigma(n)}{\sigma^*(n)} \ge \frac{\sigma^{\mathrm{e}}(n)}{\sigma^{\mathrm{e}*}(n)}, \text{ for all } n = p_1^{a_1} p_2^{a_2} \cdots p_r^{a_r} > 1, \text{ with } a_i \ge 2, \ (\forall) \ i = \overline{1, r}.$$

Case IV. Let's consider $n = n_1 n_2$, where n_1 is squarefree, and $n_2 = \prod p^a$. $p|n_2 \\ a \ge 2$

Since $(n_1, n_2) = 1$, we deduce the following relation

$$\frac{\sigma(n)}{\sigma^*(n)} = \frac{\sigma(n_1)\sigma(n_2)}{\sigma^*(n_1)\sigma^*(n_2)} \ge \frac{\sigma^{\mathrm{e}}(n_1)\sigma^{\mathrm{e}}(n_2)}{\sigma^{\mathrm{e}*}(n_1)\sigma^{\mathrm{e}*}(n_2)} = \frac{\sigma^{\mathrm{e}}(n)}{\sigma^{\mathrm{e}*}(n)}.$$
emonstration is complete.

Thus, the demonstration is complete.

Remark 2. In Theorem 2 and Theorem 3 the equality in relations (8) and (13) holds when n = 1 or n is squarefree.

Theorem 4. For any $n \ge 1$ the following inequality au

$$(n) + 1 \ge \tau^{\mathrm{e}}(n) + \tau^{*}(n)$$
 (15)

holds.

Proof. If n = 1, then we obtain $\tau(1) + 1 = 2 = \tau^{e}(1) + \tau^{*}(1)$. We consider n > 1. To prove the above inequality, will have to study several cases, namely:

Case I. If $n = p_1^2 p_2^2 \dots p_r^2$, then $\tau(n) = 3^r$ and

$$\tau^{\mathbf{e}}(n) = \tau(a_1) \cdot \tau(a_2) \cdot \dots \cdot \tau(a_r) = \tau^r(2) = 2^r = \tau^*(n),$$

which means that inequality (15) is equivalent to the inequality $3^r + 1 > 1$ $\geq 2 \cdot 2^r$, which is true because, by using *Jensen*'s inequality, we have

$$\frac{3^r + 1}{2} \ge \left(\frac{3+1}{2}\right)^r = 2^r.$$

Case II. If $a_k \neq 2$, $(\forall) k = \overline{1, r}$, then the numbers

$$\frac{n}{p_1}, \frac{n}{p_2}, \dots, \frac{n}{p_r}, \frac{n}{p_1 p_2}, \dots, \frac{n}{p_i p_j}, \dots, \frac{n}{p_i p_j p_k}, \dots, \frac{n}{p_1 p_2 \cdots p_r}$$

are not exponential divisors of n, so they are in a total number of $2^r - 1$, such that we have the inequality

$$\tau(n) = \sum_{d|_{e}n} 1 + \sum_{d\nmid_{e}n} 1 = \tau^{e}(n) + \sum_{d\nmid_{e}n} 1 \ge \tau^{e}(n) + 2^{r} - 1.$$

Therefore, we have

$$\tau(n) \ge \tau^{\mathbf{e}}(n) + 2^r - 1$$
, so $\tau(n) + 1 \ge \tau^{\mathbf{e}}(n) + \tau^*(n)$.

Case III. If there is at least one $a_j = 2$ and at least one $a_k \neq 2$, where $j, k \in \{1, 2, ..., r\}$, then without loss of the generality we renumber the prime factors from the factorization of n and we obtain

$$n = p_1^2 p_2^2 \cdots p_s^2 p_{s+1}^{a_{s+1}} \cdots p_r^{a_r}$$
, with $a_{s+1}, a_{s+2}, \dots, a_r \neq 2$

Therefore, we write $n = n_1 \cdot n_2$, where $n_1 = p_1^2 p_2^2 \cdots p_s^2$ and $n_2 = p_{s+1}^{a_{s+1}} \cdots p_r^{a_r}$, which means that $(n_1, n_2) = 1$, and

$$\begin{aligned} \tau(n) &= \tau(n_1 \cdot n_2) = \tau(n_1) \cdot \tau(n_2) \ge (\tau^{\rm e}(n_1) + \tau^*(n_1) - 1)(\tau^{\rm e}(n_2) + \tau^*(n_2) - 1) = \\ &= (\tau^{\rm e}(n_1) + 2^s - 1)(\tau^{\rm e}(n_2) + 2^{r-s} - 1) = \\ &= \tau^{\rm e}(n_1)\tau^{\rm e}(n_2) + \tau^{\rm e}(n_1)(2^{r-s} - 1) + \tau^{\rm e}(n_2)(2^s - 1) + (2^s - 1)(2^{r-s} - 1) \ge \\ &\ge \tau^{\rm e}(n) + 2^{r-s} - 1 + 2^s - 1 + (2^s - 1)(2^{r-s} - 1) \ge \tau^{\rm e}(n) + 2^r - 1 = \tau^{\rm e}(n) + \tau^*(n) - 1. \\ &\text{Thus, the inequality of the statement is true.} \end{aligned}$$

Lemma 5. For any $x_i > 0$ with $i \in \{1, 2, ..., n\}$, there is the following inequality:

$$\prod_{i=1}^{n} (1+x_i+x_i^2) + \prod_{i=1}^{n} x_i^2 \ge \prod_{i=1}^{n} (x_i+x_i^2) + \prod_{i=1}^{n} (1+x_i^2).$$
(16)

Proof. We apply the principle of mathematical induction.

Theorem 6. For any $n \ge 1$, the following inequality:

$$\sigma(n) + n \ge \sigma^{e}(n) + \sigma^{*}(n) \tag{17}$$

holds.

Proof. If n = 1, then we obtain $\sigma(1) + 1 = 2 = \sigma^{e}(1) + \sigma^{*}(1)$. Let's consider n > 1. To prove the above inequality, we will have to study several cases namely:

Case I. If $n = p_1^2 p_2^2 \cdots p_r^2$, then we deduce the equalities $\sigma(n) = \prod_{i=1}^r (1 + p_i + p_i)$

 $+p_i^2$, $\sigma^{e}(n) = \prod_{i=1}^r (p_i + p_i^2)$ and $\sigma^*(n) = \prod_{i=1}^r (1 + p_i^2)$, which means that inequality (17) implies the inequality

$$\prod_{i=1}^{r} (1+p_i+p_i^2) + \prod_{i=1}^{r} p_i^2 \ge \prod_{i=1}^{r} (p_i+p_i^2) + \prod_{i=1}^{r} (1+p_i^2),$$

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which is true, because we use inequality (16), for n = r and $x_i = p_i$. Case II. If $a_k \neq 2$, $(\forall) \ k = \overline{1, r}$, then the numbers

$$\frac{n}{p_1}, \frac{n}{p_2}, \dots, \frac{n}{p_r}, \frac{n}{p_1 p_2}, \dots, \frac{n}{p_i p_j}, \dots, \frac{n}{p_i p_j p_k}, \dots, \frac{n}{p_1 p_2 \cdots p_r}$$

are not exponential divisors of n, so they are in a total number of $2^r - 1$, and their sum is $\psi(n) - n$, so that as we have the inequality

$$\sigma(n) = \sum_{d|e^n} d + \sum_{d\nmid e^n} d = \sigma^e(n) + \sum_{d\restriction e^n} d \ge \sigma^e(n) + \psi(n) - n$$

Since we have the inequality

$$\psi(n) = n \prod_{i=1}^{r} \left(1 + \frac{1}{p_i} \right) \ge n \prod_{i=1}^{r} \left(1 + \frac{1}{p_i^{a_i}} \right) = \sigma^*(n),$$

it follows that

$$\sigma(n) \ge \sigma^{\mathbf{e}}(n) + \sigma^{*}(n) - n$$
, so $\sigma(n) + n \ge \sigma^{\mathbf{e}}(n) + \sigma^{*}(n)$.

Case III. If there is at least one $a_k \neq 2$, and at last one $a_j = 2$, where $j, k \in \{1, 2, ..., r\}$, then without loss of generality, we renumber the prime factors from the factorization of n and we obtain

$$n = p_1^2 p_2^2 \cdots p_s^2 p_{s+1}^2 \cdots p_r^{a_r}$$
, with $a_{s+1}, a_{s+2}, \dots, a_r \neq 2$.

Hence, we will write $n = n_1 \cdot n_2$, where $n_1 = p_1^2 p_2^2 \cdots p_s^2$ and $n_2 = p_{s+1}^{a_{s+1}} \cdots p_r^{a_r}$, which means that $(n_1, n_2) = 1$, and by simple calculations, it is easy to see that

$$\begin{aligned} \sigma(n) &= \sigma(n_1 \cdot n_2) = \sigma(n_1) \cdot \sigma(n_2) \ge (\sigma^{e}(n_1) + \sigma^{*}(n_1) - n_1)(\sigma^{e}(n_2) + \sigma^{*}(n_2) - n_2) = \\ &= \sigma^{e}(n_1)\sigma^{e}(n_2) + \sigma^{e}(n_1)(\sigma^{*}(n_2) - n_2) + \sigma^{*}(n_1)(\sigma^{e}(n_2) - n_2) + \sigma^{*}(n_1) \cdot \sigma^{*}(n_2) - \\ &- n_1(\sigma^{e}(n_2) + \sigma^{*}(n_2)) + n_1 n_2 \ge \sigma^{e}(n) + n_1(\sigma^{*}(n_2) - n_2) + n_1(\sigma^{e}(n_2) - n_2) + \\ &+ \sigma^{*}(n) - n_1\sigma^{e}(n_2) - n_1\sigma^{*}(n_2) + n_1 n_2 \ge \sigma^{e}(n) + \sigma^{*}(n) - n. \end{aligned}$$

Thus, the demonstration is complete.

Remark 3. In Theorem 4 and Theorem 6 the equality in relations (15) and (17) hold, when n = 1 or n is squarefree.

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Harmonic quadrilaterals revisited

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Abstract. With a plethora of instruments ranging from harmonic divisions to various geometric transformations, synthetic projective geometry has become very popular in olympiad geometry problems nowadays. In this note we expand a bit the literature concerned with applications of these techniques by discussing a topic of great importance (or better said, spectacularity) which might not be that covered in usual Euclidean geometry textbooks: harmonic quadrilaterals.

Keywords: harmonic, harmonic division, harmonic quadrilateral, inversion, symmedian.

MSC: 51A05, 51A20, 51A45.

1. Introduction

The name of **harmonic quadrilateral** dates from the middle of the nineteenth century, belonging to the famous mathematician R. Tucker [1]. He gives a slightly different definition from the ones popular these days:

Let A, B, C, D, M be five points in plane such that A, B, C, D are concyclic and M is the midpoint of the segment BD. Denote BD = 2m, i.e. MB = MD = m. Call vectorial inversion and denote (M, m^2, BD) the transformation which firstly inverts with respect to the pole M with constant m^2 , and secondly takes the image of the inversed object in BD. Then, the quadrilateral ABCD is **harmonic** if and only if each of A, C could be obtained of the other after a vectorial inversion (M, m^2, BD) .

Remark. In the above statement it isn't necessary to firstly point out the concyclity of A, B, C, D. In fact, if we perform a vectorial inversion

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 (M, m^2, BD) of a point T in plane, then its image under the transformation, T^* , lies on the circumcircle of $\triangle TBD$.

In the following, we will firstly present the most used characterization these days, a more synthetic view:

Consider ABCD a cyclic quadrilateral inscribed in $\mathcal{C}(O)$ and X a point on C. Let an arbitrary line d intersect the lines XA, XB, XC, XD at M, N, P and Q, respectively. If $\frac{NM}{NP} = \frac{QM}{QP}$, then the quadrilateral ABCDis called **harmonic**. In addition, ABCD is a harmonic quadrilateral if and only if $AB \cdot CD = BC \cdot DA$.



Proof. Denote by T the intersection of AC and BD. Since the pencil X(A, B, C, D) is harmonic, it is clear that the pencil A(A, B, C, D) is harmonic. Hence, if by S we note the intersection of BD with the tangent in A w.r.t. C, then the quadruple (SDTB) is harmonic.

Similarly, the division (S'TDB) is harmonic, where by S' we denoted the intersection of BD with the tangent in C to C. Therefore, $S \equiv S'$, i.e. BD, the tangent in A, respectively the tangent in C w.r.t. C are concurrent. But since $\triangle SAD = \triangle SBA$, if follows that SA/SB = AD/BA.

Likewise, $\triangle SCD = \triangle SBC$, i.e. SD/SC = CD/BC. Hence, $AB \cdot CD = BC \cdot DA$.

2. Another classical property

Consider ABCD a cyclic quadrilateral inscribed in C, having O as the intersection of its diagonals. Then, BO is the *B*-symmedian in $\triangle ABC$ and DO is the *D*-symmedian in $\triangle ADC$ if and only if ABCD is harmonic.



Proof. As above, one can notice that a cyclic quadrilateral is harmonic if and only the tangents in two opposite points w.r.t. the circumscribed circle concur on the diagonal determined by the other two points, i.e. BB, DD and AC are concurrent in a point X (where by BB we denote the tangent in B to the circumcircle of ABC).

On the other hand, it is well-known that in a triangle MNP, having M_1 as the foot of the M-symmedian and M_2 as the intersection of the tangent in M to its circumcircle with the side BC, the quadruple (M_2, N, M_1, P) is a harmonic division. Hence, BO is the B-symmedian in $\triangle ABC$ and DO is the D-symmedian in $\triangle ADC$ if and only if BB, DD and AC are concurrent.

Remark. In the initial statement it isn't really necessary to stress from the beginning that ABCD is cyclic. One can give the following enuntiation:

Consider ABCD a convex quadrilateral inscribed in C, having O as the intersection of its diagonals. Then, BO is the B-symmedian in $\triangle ABC$ and DO is the D-symmedian in $\triangle ADC$ if and only if ABCD is harmonic.

As probably yet noticed, the condition of BO, DO being the symmetians in $\triangle ABC$, respectively $\triangle ADC$, together with the harmonicity of (X, A, O, C), imply the coincidence of the polars of X w.r.t. $\triangle ABC$ and $\triangle ADC$ (the common polar beeing the line BD). Hence, the two circles coincide as well.

Corollary 1. Consider ABC a triangle and T the intersection of the tangents at B and C w.r.t. the circumcircle of ABC. Then AT is the A-symmedian.

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Corollary 2. In the initially described configuration, AO, CO are the A, respectively C-symmetry in $\triangle BAD$, $\triangle BCD$.

3. An IMO-type alternative definition

Consider ABC a triangle and D a point on its circumcircle. Draw the Simson line of D w.r.t. $\triangle ABC$. This line cuts its sides BC, CA, AB in P, Q and R, respectively. Then, the quadrilateral ABCD is **harmonic** if and only if PQ = QR.



The above statement is slightly modifying the original one from IMO 2003, problem 4:

Let ABCD be a cyclic quadrilateral. Let P, Q and R be the feet of the perpendiculars from D to the lines BC, CA and AB, respectively. Show that PQ = QR if and only if the bisectors of $\triangleleft ABC$ and $\triangleleft ADC$ meet on AC.

First proof. As we first stated, it is well-known that P, Q, R are collinear on the Simson line of D. Moreover, since $\measuredangle DPC$ and $\measuredangle DQC$ are right angles, the points D, P, Q, C are concyclic and so $\measuredangle DCA = \measuredangle DPQ = \measuredangle DPR$. Similarly, since D, Q, R, A are concyclic, we have $\measuredangle DAC = \measuredangle DRP$. Therefore, $\triangle DCA \sim \triangle DPR$.

Likewise, $\triangle DAB \sim \triangle DQP$ and $\triangle DBC \sim \triangle DRQ$. Then

$$\frac{DA}{DC} = \frac{DR}{DP} = \frac{DB \cdot \frac{QR}{BC}}{DB \cdot \frac{PQ}{BA}} = \frac{QR}{PQ} \cdot \frac{BA}{BC}.$$

0 D

Thus PQ = QR if and only if DA/DC = BA/BC, whence by the converse of the bisector theorem, we deduce that it is equivalent with the concurrence of the bisectors of $\measuredangle ABC$ and $\measuredangle ADC$ on AC.

Second proof. Because $DP \perp BC$, $DQ \perp AC$, $DR \perp AB$, the circles with diameters DC and DA contain the pairs of points P, Q and Q, R, respectively. It follows that $\measuredangle PDQ$ is equal to γ or $180^{\circ} - \gamma$, where $\gamma = \measuredangle ACB$.

Similarly, $\measuredangle QDR$ is equal to α or $180^{\circ} - \alpha$, where $\alpha = \measuredangle CAB$. Then, by the law of sines, we have $PQ = CD \sin \gamma$ and $QR = AD \sin \alpha$. Hence the condition PQ = QR is equivalent with $CD/AD = \sin \alpha / \sin \gamma$.

On the other hand, $\sin \alpha / \sin \gamma = CB/AB$ by the law of sines again,. Thus PQ = QR if and only if CD/AD = CB/AB, i.e. $AB \cdot CD = CB \cdot AD$.

4. A probably new non-standard characterization

Let $\mathcal{C}(\mathcal{O})$ be a given circle and X, Y two points in its plane. It is wellknown that there exist exactly two circles, ρ_1 and ρ_2 , passing through X, Ytangent to the initial circle \mathcal{C} . Consider A_1 and A_2 the mentioned tangency points. Draw a third mobile circle, ρ_3 , through X, Y touching the reference circle \mathcal{C} at B_1 and B_2 , respectively. Then, the quadrilateral $A_1B_1A_2B_2$ is harmonic.



Proof. For example, perform an inversion with pole X. Hence, C is mapped into an other circle C^* and ρ_1 , ρ_2 into the lines ρ_1^* and ρ_2^* , respectively, not passing through X, representing the tangents from Y^* to C^* , i.e. A_1^* , A_2^* are the tangency points of C^* with ρ_1^* , ρ_2^* .

Since ρ_3 is mapped into the line ρ_3^* through Y^* cutting \mathcal{C}^* at B_1^*, B_2^* , we deduce that the harmonicity of $A_1B_1A_2B_2$ is equivalent with the harmonicity of $A_1^*B_1^*A_2^*B_2^*$, which is clear, because $A_1^*A_2^*$ is the polar of Y^* w.r.t. \mathcal{C}^* . \Box

If in the above configuration we consider \mathcal{C} a circle centered at infinity, we deduce the following consequence:

Corollary. Consider d a line and X, Y two points in its plane. Denote by A_1 , A_2 the points where the two circles passing through X, Y and tangent to d, touch the respective line. Let ω be an arbitrary circle through X, Y, which touches d at B_1 and B_2 . Then, the quadruple (A_1, B_1, A_2, B_2) forms a harmonic division.



Their converses are also true. Precisely, if we consider A_1 , A_2 similarly as above and this time B_1 , B_2 are two points on C such that $A_1B_1A_2B_2$ is a harmonic quadrilateral (degenerated or not), then the points X, Y, B_1 , B_2 lie on a same circle.

5. The Brocard points of a harmonic quadrilateral

In this section we present the existence of the Brocard points of a harmonic quadrilateral, proved by F. G. W. Brown [2], more as an other characterization of this particular cyclic quadrilateral.



Consider ABCD a cyclic quadrilateral. If there exists a point X such that the angles $\triangleleft XAD$, $\triangleleft XBA$, $\triangleleft XCB$, $\triangleleft XDC$ are equal, then quadrilateral ABCD is harmonic.

Proof. Denote the sides BC, CD, DA, AB by a, b, c, d; the diagonals DB, AC by e, f; the area of the quadrilateral by Q; the circumradius by R and each of the angles XAD, XBA, XCB, XDC by ω . Then

$$\sphericalangle AXB = \pi - \omega - (A - \omega) = \pi - A$$

Similarly, $\sphericalangle BXC = \pi - B$, $\measuredangle CXD = \pi - C$, $\measuredangle DXA = \pi - D$. Now, since

$$\frac{AX}{\sin\omega} = \frac{d}{\sin AXB} = \frac{d}{\sin(\pi - A)} = \frac{d}{\sin A}$$

and

$$\frac{AX}{\sin(D-\omega)} = \frac{c}{\sin AXD} = \frac{c}{\sin(\pi-D)} = \frac{c}{\sin D}$$

we deduce that

$$\frac{\sin\left(D-\omega\right)}{\sin\omega} = \frac{d\sin D}{c\sin A},$$

or

$$\cot \omega = \frac{d}{c \sin A} + \cot D.$$

Similarly,

$$\cot w = \frac{c}{b \sin D} + \cot C, \qquad (*)$$

and the cyclic analogues.

Hence,

$$\cot w = \frac{d}{c \sin A} + \cot D = \frac{c}{b \sin D} + \cot C = \frac{b}{a \sin C} + \cot B = \frac{a}{d \sin B} + \cot A$$
or

$$\frac{d}{c\sin A} - \frac{a}{d\sin B} + \frac{b}{a\sin C} - \frac{c}{b\sin D} = \cot A - \cot B + \cot C - \cot D$$

Since ABCD is cyclic, $1/\sin A = 1/\sin C$, $\cot A = -\cot C$, etc; therefore

$$\left(\frac{d}{c} + \frac{b}{a}\right)\frac{1}{\sin A} = \left(\frac{c}{b} + \frac{a}{d}\right)\frac{1}{\sin B}$$

But

$$(ab+cd)\sin A = (ad+bc)\sin B = 2Q;$$

hence ac = bd, i.e. ABCD is harmonic.

Its converse is also valid and can be proved on the same idea. Moreover, there is a second Brocard point X', such that $\sphericalangle X'AB = \sphericalangle X'BC =$

 $= \measuredangle X'CD = \measuredangle X'DA = \omega'$. By simple manipulations on (*) and using *Ptolemy*'s theorem, we obtain

$$\cot \omega = \frac{8R^2Q}{e^2f^2}.$$

Similarly, we deduce that

$$\cot \omega' = \frac{8R^2Q}{e^2f^2}$$

Hence, $\omega = \omega'$.

Further, we will compute the distance XX'. For this, let P be the foot of the perpendicular from X to BC and CP = x, PX = y. Then

$$x = CX\cos\omega = \frac{b\sin\omega\cos\omega}{\sin C} = \frac{b(ab+cd)\cot\omega}{2Q(1+\cot^2\omega)} = \frac{bR}{c} \cdot \sin 2\omega$$

and

$$y = x \tan \omega = \frac{b(ab + cd)}{2Q(1 + \cot^2 \omega)} = \frac{2bR}{c} \sin^2 \omega$$

Similarly, if x' = P'B, y' = P'X', where P' is the foot of the perpendicular from X' to BC, then

$$x' = \frac{d(ad+bc)\cot\omega}{2Q(1+\cot^2\omega)}$$

and

$$y' = \frac{d(ad+bc)}{2Q(1+\cot^2\omega)}$$

Since by the law of cosines $XX'^2 = XB^2 + X'B^2 - 2XB\cdot X'B\cdot\cos XBX'$ and by

$$\sin B = \frac{ab + cd}{2bd} \tan \omega = \frac{a^2 + d^2}{2ad} \cdot \tan \omega$$

using that ac = bd, i.e. ef = 2ac = 2bd, it follows that

$$XX' = \frac{efbd\cos\omega \cdot \sqrt{\cos 2\omega}}{2RQ} = \frac{b^2d^2\cos\omega \cdot \sqrt{\cos 2\omega}}{RQ}.$$

6. Remarks on an Iranian concurrence problem

This problem firstly became popular on the Mathlinks forum for its difficulty. It is a nice result involving a concurrence of three lines in triangle.

Let ABC be a triangle. The incircle of $\triangle ABC$ touches the side BC at A', and the line AA' meets the incircle again at a point P. Let the lines CP and BP meet the incircle of triangle ABC again at N and M, respectively. Prove that the lines AA', BN and CM are concurrent.

Iran National Olympiad 2002



Proof. Consider Q the intersection of MN with BC. By Ceva and Menelaus theorem, it is clear that PA', BN and CM concur if and only if the quadruple (Q, B, A', C) is harmonic. So, the problem reduces to proving that $Q \equiv Q'$, where Q' is the intersection of the tangent in P to the incircle with BC.

Since (Q', B, D, C) is a harmonic division, we deduce that the pencil P(Q', B, D, C) is harmonic and by intersecting it with the incircle, it follows that the quadrilateral PMA'N is harmonic. Hence, the lines MN, the tangent in P, respectively, the tangent in A' to the incircle are concurrent, i.e. $Q \equiv Q'$.

We leave the readers a more difficult extended result (Figure 8.):

I. Suppose that in the configuration described above the lines AA', BN and CM are concurrent in a point X. Similarly, one can prove the excircle related problem, meaning, the lines AA'', BN_1 , CM_1 are concurrent in a point X_1 , where by A'' we denoted the tangency point of the A-excircle with the side BC and by N_1 , M_1 the second intersections of CP', respectively BP' with the A-excircle, where P' is the second intersection of the line AA'' with the A-excircle. Prove that the lines PP', BC and XY are concurrent.

Moreover, we can go further and observe the next concurrence, which remains as an other proposed problem for those who are interested (Figure 9.):

II. Consider that in the above configuration the lines PP', BC and XY concur in a point X_a . Similarly, we deduce the existence of points Y_b , Z_c , with same right as X_a , on the lines CA, respectively AB. Then, the triangles ABC and $X_aY_bZ_c$ are perspective (i.e. the lines AX_a , BY_b , CZ_c concur); their perspector is X(75), i.e. the isotomic conjugate of the incenter.

The last remark, regarding the determination of the perspector, was communicated by Darij Grinberg, using computer dynamic geometry.



Figure 9.

In the following we leave again the readers, this time, a slightly more general statement of the original problem:

III. Let ABC be a triangle. The incircle of $\triangle ABC$ touches the side BC at A', and let P be an arbitrary point on (AA'). Let the lines CP and BP meet the incircle of triangle ABC again at N, N' and M, M', respectively. Prove that the lines AA', BN and CM, respectively AA', BN', CM', are concurrent. (i.e. the lines BC, MN, M'N', B'C' are concurrent, where B', C' are the tangency points of the incircle with CA, respectively AB.)



Figure 10

As a remark, it is clear that an extension similar to **II** of the above result is usually **not** true.

7. On a perpendicularity deduced projectively

This problem is a result involving a perpendicularity as a consequence of a projective fact. Moreover, we will give a pure by synthetic proof, based on a simple angle chasing.



Let $\rho(O)$ be a circle and A a point outside it. Denote by B, C the points where the tangents from A w.r.t to $\rho(O)$ meet the circle, D the point

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on $\rho(O)$ for which $O \in (AD)$, X the foot of the perpendicular from B to CD, Y the midpoint of the line segment BX and by Z the second intersection of DY with $\rho(O)$. Prove that $ZA \perp ZC$ (Figure 11).

BMO Shortlist 2007, proposed by the author of this paper

First proof. Let H the second intersection of CO with $\rho(O)$. Thus $DC \perp DH$, so $DH \parallel BX$.

Because Y is the midpoint of (BX), we deduce that the division (B, Y, X, ∞) is harmonic, so also is the pencil D(B, Y, X, H), and by intersecting it with $\rho(O)$, it follows that the quadrilateral HBZC is harmonic.

Hence the pencil C(HBZC) is harmonic, so by intersecting it with the line HZ, it follows that the division (A'ZTH) is harmonic, where A', T are the intersections of HZ with the tangent in C, respectively with BC.

So, the line CH is the polar of A' w.r.t. $\rho(O)$, but $CH \equiv BC$ is the polar of A as well, so $A \equiv A'$, hence the points H, Z, A are collinear, therefore $ZA \perp ZC$.

Second proof. Let E be the midpoint of BC. So $EY \parallel CD$ and $EY \perp BX$. Since $\triangleleft YEB = \triangleleft DCB = \triangleleft DZB$, we deduce that Y, E, Z, B lie on a same circle. Thus $\triangleleft ZBC = \triangleleft ZEA$ (OA is tangent to the circle YEZB).

But because $\sphericalangle ZBC = \sphericalangle ZCA$, we have that $\sphericalangle ZCA = \sphericalangle ZEA$, i.e. the points C, A, Z, E are concyclic. Hence $\sphericalangle CZA = \sphericalangle CEA = 90$, i.e. $ZA \perp ZC$.

On the idea of the second solution, we give a more general statement:

IV. The tangents from a point A to the circle $\rho(O)$ touch it at the points B, C. For a point $E \in (BC)$ denote: the point D on ρ for which $E \in (AD)$; the second intersection X of the line CD with the circumcircle of the triangle BED; the point Y on the line BX for which EY || CD; the second intersection Z of DY with the circle ρ ; the intersection T of the line AC with BX; the second intersection H of the line AZ with the circle ρ . Prove that the quadrilaterals AZEC, AZYT are cyclic and BH || \overline{AED} . (Figure 12)

Proof. Since $\triangleleft ZBE = \triangleleft ZBC = \triangleleft ZDC = \triangleleft ZYE$, we have that $\triangleleft ZBE = \triangleleft ZYE$, i.e. the quadrilateral ZBYE is cyclic. Hence $\triangleleft EBY = \triangleleft EBX = \triangleleft EDX = \triangleleft EDC = \triangleleft YED$. So $\triangleleft EBY = \triangleleft YED$, meaning that the line \overline{AED} is tangent to the circumcircle of the quadrilateral ZBYE.

Thus, $\triangleleft ZEA = \triangleleft ZBE = \triangleleft ZBC = \triangleleft ZCA$, i.e. the quadrilateral AZEC is cyclic. So, $\triangleleft ZYB = \triangleleft ZEB = \triangleleft ZAC = \triangleleft ZAT$, whence the quadrilateral AZYT is also cyclic.

Finally, we have $\measuredangle AEC = \measuredangle AZC = \measuredangle CDH$. Hence, $BH \| \overline{AED}$. \Box



Observe that when E is the midpoint of the segment BC we obtain the original proposed problem.

8. On a classic locus problem as a recent IMO Team Preparation Test exercise

Consider the isosceles triangle ABC with AB = AC, and M the midpoint of BC. Find the locus of the points P interior to the triangle, for which $\sphericalangle BPM + \sphericalangle CPA = \pi.$

IMAR Contest 2006



Figure 13

First proof (**by Dan Schwarz**). We will start with the following claim:

Lemma. Prove that if, for a point P interior to the triangle, $\triangleleft ABP = \triangleleft BCP$, then $\triangleleft BPM + \triangleleft CPA = \pi$.

Proof. Notice that the configuration is fully symmetrical:

and also that $P \in AM$ guarantees the result, therefore we may assume w.l.o.g. (for construction's sake) that P is interior to the triangle ABM. The given angle equality is readily seen to be equivalent with $P \in \Gamma$, where Γ is the arc interior to ABC of the circumcircle \mathcal{K} of triangle BCI, with I the incenter of ABC. It is also immediately seen that AB and AC are tangent to \mathcal{K} .

Consider the Apollonius circle \mathcal{A} for points A, M and ratio BA/BM; denote by U and V the points where AM meets \mathcal{A} . Clearly, $B, C \in \mathcal{A}$. Then BA/BM = UA/UM = VA/VM, from which follows $BA^2/BM^2 =$ $= UA \cdot VA/UM \cdot VM = UA \cdot VA/BM^2$, by symmetry and the power of a point relation, so $BA^2 = UA \cdot VA$, hence AB is tangent to \mathcal{A} , again by the power of a point relation. Therefore the two circles \mathcal{K} and \mathcal{A} coincide. Now, prolong AP until it meets \mathcal{K} again at Q, and prolong PM until it meets \mathcal{K} again at R. Then BA/BM = PA/PM = QA/QM, therefore $BA^2/BM^2 = PA \cdot QA/PM \cdot QM = BA^2/PM \cdot QM$, by the power of a point relation, so $BM^2 = PM \cdot QM$. But $BM^2 = PM \cdot RM$, again by the power of a point relation, so QM = RM, hence BQ = CR. It follows that $\sphericalangle BPQ = \measuredangle CPR$, and this is enough to yield $\measuredangle BPM + \measuredangle CPA = \pi$.

Alternatively, one can calculate ratios using the symmetrical points P'and Q' with respect to AM, and denoting by N the meeting point of PQ'and P'Q. It can be obtained that AM = AN, therefore $M \equiv N$ and the rest easily follows as above.

Returning to the problem, we claim the locus is the arc Γ (defined in the above), together with the (open) segment (AM). Clearly $P \in (AM)$ fills the bill, so from now on we will assume $P \notin (AM)$. Also, as above, we will assume P interior to the triangle ABM (otherwise we work with the symmetrical relations).

Assume $P \notin \Gamma$, equivalent to $\measuredangle BPC \neq \pi - \measuredangle ABC$; then AB and ACare not tangent to \mathcal{K} . Take B' and C' to be the tangency points on \mathcal{K} from A, and M' to be the midpoint of B'C'. We are now under the conditions from the lemma (for triangle AB'C'), therefore $\measuredangle B'PM' + \measuredangle C'PA = \pi$. But $\measuredangle B'PM' = \measuredangle BPM + \delta(\measuredangle B'PB + \measuredangle MPM')$ and $\measuredangle C'PA = \measuredangle CPA - -\delta \measuredangle C'PC$, where $\delta = 1$ if $\measuredangle BPC < \pi - \measuredangle ABC$, respectively $\delta = -1$ if $\measuredangle BPC > \pi - \measuredangle ABC$. We have $\measuredangle B'PB = \measuredangle C'PC$ from the symmetry of the configuration, and $\measuredangle BPM + \measuredangle CPA = \pi$ given; these relations therefore imply $\triangleleft MPM' = 0$, which can only happen when M, M' and P are collinear, i.e. $P \in AM$, which was ruled out from the start in this part of the proof. The contradiction thus reached confirms our claim.

Second proof. Denote the point D as the intersection of the line AP with the circumcircle of BPC and $S = DP \cap BC$.

Since $\triangleleft SPC = 180^{\circ} - \triangleleft CPA$, it follows that $\triangleleft BPS = \triangleleft CPM$.

From the *Steiner* theorem applied to triangle BPC for the isogonals PS and PM,

$$\frac{SB}{SC} = \frac{PB^2}{PC^2}.$$

On other hand, using the law of sines, we obtain

$$\frac{SB}{SC} = \frac{DB}{DC} \cdot \frac{\sin SDB}{\sin SDC} = \frac{DB}{DC} \cdot \frac{\sin PCB}{\sin PBC} = \frac{DB}{DC} \cdot \frac{PB}{PC}.$$

Thus by the above relations, it follows that DB/DC = PB/PC, i.e. the quadrilateral PBDC is harmonic. Therefore the point $A' = BB \cap CC$ lies on the line PD (where by XX we denoted the tangent in X to the circumcircle of BPC).

If A' = A, then lines AB and AC are always tangent to the circle BPC, and so the locus of P is the circle BIC, where I is the incenter of ABC. Otherwise, if $A' \neq A$, then $A' = AM \cap PS \cap BB \cap CC$, due to the fact that A' lies on PD and $A = PS \cap AM$, and by maintaining the condition that $A' \neq A$, we obtain that $PS \equiv AM$, therefore P lies on (AM). \Box

Next, after two synthetic approaches, we continue with a direct trigonometric solution, actually the one which all three contestants who solved the problem used.

Third solution. Consider $m(\sphericalangle BPM) = \alpha$ and $m(\sphericalangle CPM) = \beta$ so $m(\sphericalangle APC) = 180^{\circ} - \alpha$ and $m(\sphericalangle APB) = 180^{\circ} - \beta$. Also let $m(\sphericalangle CBP) = u$, $m(\sphericalangle BCP) = v$ and denote AB = AC = l.

By the law of sines, applied to triangles ABP and ACP, $\frac{l}{\sin(180 - \beta)} = AP$

$$= \frac{AP}{\sin(B-u)} \text{ and } \frac{l}{\sin(180-\alpha)} = \frac{AP}{\sin(C-v)}. \text{ Hence}$$
$$\frac{\sin\alpha}{\sin\beta} = \frac{\sin(C-v)}{\sin(B-u)}. \tag{(\star)}$$

Again, by the law of sines, this time applied in to the triangles BPMand CPM, $\frac{PM}{\sin u} = \frac{BM}{\sin \alpha}$ and $\frac{PM}{\sin v} = \frac{CM}{\sin \beta}$. Hence $\frac{\sin \alpha}{\sin \beta} = \frac{\sin u}{\sin v}$. (**) From (*) and (**) we deduce that $\frac{\sin u}{\sin v} = \frac{\sin(C-v)}{\sin(B-u)}$, i.e. $\cos(B-2u) - \cos B = \cos(C-2v) - \cos C$. Since B = C, it follows that $\cos(B-2u) = \cos(B-2v)$.

If B - 2u = B - 2v, then we have u = v i.e. P lies on (AM). Alternatively, if B - 2u = 2v - B, then we have B = u + v, i.e. P lies on the arc BIC.

9. The *orthotransversal* line of a triangle refreshed as a Mathematical Reflections problem

Let ABC be a triangle and P be an arbitrary point inside the triangle. Let A', B', C' be respectively the intersections of AP and BC, BP and CA, CP and AB. Through P we draw a line perpendicular to PA that intersects BC at A_1 . We define B_1 and C_1 analogously. Let P' be the isogonal conjugate of the point P with respect to triangle A'B'C'. Then A_1 , B_1 , C_1 all lie on a same line l that is perpendicular to PP'.

Math. Reflections, 4/2006, problem O13, proposed by Khoa Lu Nguyen

In [7], Bernard Gibert names the line $\overline{A_1B_1C_1}$ the orthotransversal line of P.



First proof. Given four collinear points X, Y, Z, T, let (X, Y, Z, T) denote, we mean the cross-ratio of four points X, Y, Z, T. Given four concurrent lines x, y, z, t, let (x, y, z, t) denote, we mean the cross-ratio of four lines x, y, z, t. We first introduce the following claim:

Lemma. Let ABC be a triangle and P' be the isogonal conjugate of an arbitrary point P with respect to ABC. Then the six projections from Pand P' to the sides of triangle ABC lie on a circle with center the midpoint of PP'.

Proof. Let P_a , P_b , P_c be the projections from P to the sides BC, CA, AB. Similarly, let P'_a , P'_b , P'_c be the projections from P' to the sides BC, CA, AB. Call O the midpoint of PP'. We need to show that P_a , P_b , P_c , P'_a , P'_b , P'_c lie on a circle with center O.

Consider the trapezoid $PP'P'_aP_a$ that has $m(PP_aP'_a) = m(P'P'_aP_a) = 90^{\circ}$ and O the midpoint of (PP'). Hence, O must lie on the perpendicular bisector of $P_aP'_a$. By a similar argument, we obtain that O also lies on the perpendicular bisector of $P_bP'_b$ and $P_cP'_c$.

Because P' is the isogonal conjugate of P with respect to ABC, we have $\langle BAP = \langle P'AC \rangle$, or $\langle P_cAP = \langle P'AP'_b \rangle$. Hence, we obtain $\langle APP_c = \langle P'_bP'A \rangle$. On the other hand, since quadrilaterals AP_cPP_b and $AP'_cP'P'_b$ are cyclic, it follows that $\langle APP_c = \langle AP_bP_c \rangle$ and $\langle P'_bP'A = \langle P'_bP'_cA \rangle = \langle P'_bP'_cP_c \rangle$. Thus, $\langle AP_bP_c = \langle P'_bP'_cP_c \rangle$. This means that $P_bP_cP'_cP'_b$ is inscribed in a circle. We notice that the center of this circle is the intersection of the perpendicular bisectors of $P_bP'_b$ and $P_cP'_c$, which is O. In the same manner, we obtain that $P_cP_aP'_cP'_a$ is inscribed in a circle with center O. Thus these two circles are congruent as they have the same center and pass through a common point P_c . Therefore, these six projections all lie on a circle with center O.

Back to the problem, let P_a , P_b , P_c be the projections from P to the sides B'C', C'A', A'B' and O be the midpoint of PP'. Then (O) is the circumcircle of triangle $P_aP_bP_c$. We will show a stronger result. In fact, l is the polar of P w.r.t. the circle (O). We will prove that A_1 lies on the polar of P w.r.t. the circle (O), and by a similar argument so does B_1 , respectively C_1 .

Let B_2 and C_2 be the intersections of the line PA_1 with A'C' and A'B', respectively. Denote by M and N the intersections of the line PA_1 with the circle (O). It is suffice to prove that $(A_1, P, M, N) = -1$.

Consider X the projection from P to BC. Then five points P, P_b , P_c , X, A' lie on the circle (a) with diameter PA'. Since $(A'A_1, A'A, A'C', A'B') =$ = -1, we obtain that PP_cXP_b is a harmonic quadrilateral. This yields that P_bP_c and the tangents at P and X of the circle (a) are concurrent at a point U. Since the tangent at P of the circle (a) is PA_1 , it follows that U is the concurrence point of P_bP_c , PA_1 and the perpendicular bisector of (PX). Consider the right triangle PXA_1 at X. Since U lies on PA_1 and the perpendicular bisector of (PX), we deduce that U is the midpoint of (PA_1) . Therefore U lies on the line P_bP_c . Since $A'P_b\Delta \cdot A'B_2 = A'P_c \cdot A'C_2 = A'P_2$, it follows that $P_bP_cC_2B_2$ is cyclic. Hence, $UB_2 \cdot UC_2 = UP_b \cdot UP_c$. On the other hand, we have $UP_b \cdot UP_c = UM \cdot UN$, because P_bP_cMN is cyclic.

Thus, we obtain $UB_2 \cdot UC_2 = UM \cdot UN$.

Again, since $(A'A_1, A', A'C', A'B') = -1$, we have $(A_1, P, B_2, C_2) = -1$. But, because U is the midpoint of A_1P , it follows that $UB_2 \cdot UC_2 = UA_1^2$. Hence $UA_1^2 = UM \cdot UN$, i.e. $(A_1, P, M, N) = -1$.

Second proof for the collinearity of A_1 , B_1 , C_1 . (by Darij Grinberg, after Jacques Hadamard [6]) Let \mathcal{K} be an arbitrary circle centered at P; the polars of the points A, B, C, A_1 , B_1 , C_1 w.r.t. \mathcal{K} are called a, b, c, a_1 , b_1 , c_1 . After the construction of polars, we have $a \perp PA$, $b \perp PB$, $c \perp PC$, $a_1 \perp PA_1$, $b_1 \perp PB_1$ and $c_1 \perp PC_1$. Since the point A_1 lies on BC, the polar a_1 passes through the intersection $b \cap c$. From $a_1 \perp PA_1$ and $PA_1 \perp PA$, we have $a_1 || PA$; from $a \perp PA$, thus we get $a_1 \perp a$. Hence, a_1 is the line passing through the intersection $b \cap c$ and orthogonal to a. Analogously, b_1 is the line passing through the intersection $a \cap b$ and orthogonal to b, and c_1 is the line passing through the intersection $a \cap b$ and orthogonal to c. Therefore, a_1 , b_1 and c_1 are the altitudes of the triangle formed by the lines a, b and c; consequently, the lines a_1 , b_1 and c_1 concur. From this, we derive that the points A_1 , B_1 and C_1 are collinear.

Following the ideas presented above, we leave the readers a similar collinearity:

V. Consider ABC a triangle and P a point in its plane. Let R, S, T, X, Y, Z be the midpoints of the segments BC, CA, AB, PA, PB and PC, respectively. Draw the lines through the X, Y, Z, orthogonal to PA, PB, respectively PC; these lines touch ST, TR, RS at A_1 , B_1 and C_1 , respectively. Prove that the points A_1 , B_1 , C_1 lie on the same line.



10. Bellavitis' theorem on balanced quadrilaterals

We call ABCD a balanced [8] quadrilateral if and only if $AB \cdot CD = BC \cdot DA$. Particulary, if ABCD is cyclic we come back to the harmonic quadrilateral. However, this class contains other quadrilaterals, for example all kites.

Let the lengths of the sides AB, BC, CD and DA of a (convex) quadrilateral ABCD be denoted by a, b, c and d respectively. Similarly, the lengths of the quadrilateral's diagonals AC and BD will be denoted by e and f. Let Ebe the point of intersection of the two diagonals. The magnitude of $\measuredangle DAB$ will be referred to as α , with similar notation for the other angles of the quadrilateral. The magnitudes of $\measuredangle DAC$, $\measuredangle ADB$ etc. will be denoted by α_B, δ_C a.s.o. (see Figure 16.). Finally, the magnitude of $\measuredangle CED$ will be referred to as ϵ .



Figure 16.

Theorem. (Bellavitis 1854) If a (convex) quadrilateral *ABCD* is *balanced*, then

$$\alpha_B + \beta_C + \gamma_D + \delta_A = \beta_A + \gamma_B + \delta_C + \alpha_D = 180^\circ.$$

Note that the convexity condition is a necessary one. The second equality sign does not hold for non-convex quadrilaterals.

Proof. (by Eisso J. Atzema [8]) Although in literature it is known that Giusto Bellavitis himself gave a proof using complex numbers, a trigonometric proof of his theorem follows from the observation that by the law of sines for any balanced quadrilateral we have

$$\sin \gamma_B \cdot \sin \alpha_D = \sin \alpha_B = \sin \gamma_D,$$

or

$$\cos\left(\gamma_B + \alpha_D\right) - \cos\left(\gamma_B - \alpha_D\right) = \cos\left(\alpha_B + \gamma_D\right) - \cos\left(\alpha_B - \gamma_D\right)$$

That is,

$$\cos(\gamma_B + \alpha_D) - \cos(\gamma_B - \alpha + \alpha_B) = \cos(\alpha_B + \gamma_D) - \cos(\alpha_B - \gamma + \gamma_B),$$

 $\cos(\gamma_B + \alpha_D) + \cos(\delta + \alpha) = \cos(\alpha_B + \gamma_D) + \cos(\delta + \gamma).$ By cycling through, we also have

 $\cos (\delta_C + \beta A) + \cos (\alpha + \beta) = \cos (\beta_C + \delta_A) + \cos (\alpha + \delta).$ Since $\cos (\alpha + \beta) = \cos (\delta + \gamma)$, adding these two equations gives $\cos (\gamma_B + \alpha_D) + \cos (\delta_C + \beta_A) = \cos (\alpha_B + \gamma_D) + \cos (\beta_C + \delta_A),$

or

$$\cos \frac{1}{2} (\delta_C + \gamma_B + \beta_A + \alpha_D) \cdot \cos \frac{1}{2} (\gamma_B + \alpha_D - \delta_C - \beta_A)$$
$$= \cos \frac{1}{2} (\alpha_B + \beta_C + \gamma_D + \delta_A) \cdot \cos \frac{1}{2} (\alpha_B + \gamma_D - \beta_C - \delta_A).$$

Now, note that

$$\gamma_B + \alpha_D - \delta_C - \beta_A = 360^\circ - 2\epsilon - \delta - \beta$$

and likewise

$$\alpha_B + \gamma_D - \beta_C - \delta_A = 2\epsilon - \beta - \delta.$$

Finally,

$$\frac{1}{2}\left(\delta_C + \gamma_B + \beta_A + \alpha_D\right) + \frac{1}{2}(\alpha_B + \beta_C + \gamma_D + \delta_A) = 180^\circ.$$

It follows that

$$\cos\frac{1}{2}(\delta_C + \gamma_B + \beta_A + \alpha_D) \cdot \cos\left(\epsilon + \frac{1}{2}(\beta + \delta)\right) =$$
$$= -\cos\frac{1}{2}(\delta_C + \gamma_B + \beta_A + \alpha_D) \cdot \cos\left(\epsilon - \frac{1}{2}(\beta + \delta)\right),$$
$$\cos\frac{1}{2}(\delta_C + \gamma_B + \beta_A + \alpha_D) \cdot \cos\left(\epsilon\right)\cos\frac{1}{2}(\delta + \beta) = 0.$$

or

Clearly, if neither of the last two factors is equal to zero, the first factor has to be zero and we are done. The last factor, however, will be zero if and only if ABCD is cyclic. It is easy to see that any such quadrilateral has the angle property of *Bellavitis*' theorem. Therefore, in the case that ABCD is cyclic, *Bellavitis*' theorem is true. Consequently, we may assume that ABCDis not cyclic and that the third term does not vanish. Likewise, the second factor only vanishes in case ABCD is orthodiagonal. For such quadrilaterals, we know that $a^2 + c^2 = b^2 + d^2$. In combination with the initial condition ac = bd, this implies that each side has to be congruent to an adjacent side. In other words, ABCD has to be a kite. Again, it is easy to see that in that case *Bellavitis*' theorem is true. We can safely assume that ABCD is not a kite and that the second term does not vanish either. Moreover, for a synthetic solution one can see Nikolaos Dergiades' proof in [9]. Also, a solution using inversion exists. See [11].

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Problems and Solutions from SEEMOUS 2011 Competition

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Abstract. In this note, we present the problems with solutions and comments from the 5th South Eastern European Mathematical Olympiad for University Students, SEEMOUS 2011, organized by the MASEE and Romanian Mathematical Society between March 2 and March 6, 2011.

Keywords: Chebyshef's integral inequality, eigenvalues of a matrix, characteristic polynomial, vectors, Lagrange multipliers, Riemann sum, Taylor expansion.

MSC: 15A24, 15A26, 26A42, 26D15, 51D20.

1. Introducere

Cea de-a cincea ediție a Olimpiadei de Matematică pentru Studenți din Sud-Estul Europei, SEEMOUS 2011, a fost organizată de Mathematical Society of South Eastern Europe (MASSEE) și de Societatea de Științe Matematice din Romania (S.S.M.R.). Gazda competiției a fost Universitatea Politehnica din Bucuresti (U.P.B.). Trebuie mentionat sprijinul pe care conducerea acestei universități îl acordă în mod constant și eficient concursurilor studențești de matematică. În acest sens, reamintim faptul că reluarea Concursului de matematică pentru studenți "Traian Lalescu" acum patru ani s-a datorat cooperării dintre Ministerul Educației și Cercetării, Societatea de Stiinte Matematice din România și Universitatea Politehnica din București. Organizatorii SEEMOUS 2011 au beneficiat și de ajutorul colegilor din Facultatea de Matematică și Informatică din Universitatea București și al celor din Institutul de Matematică "Simion Stoilow" al Academiei Romane. Majoritatea participanților au fost de acord că această ediție a fost cea mai bine organizată, atât din punct de vedere stiințific cât și din punct de vedere logistic.

Concursul propriu-zis s-a bucurat de un număr record de participanți. Au participat 85 de studenți organizați în 20 de echipe de la universități din țări precum Bulgaria, Columbia, Grecia, Macedonia, Moldova, România, Rusia și Ucraina.

Studenții români au avut din nou o comportarea remarcabilă reușind să obțină 2 medalii de aur prin *Pădureanu Victor* (locul 2 după punctaj) (Academia Tehnică Militară, București) și *Cocalea Andrei* (Universitatea din

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Bucuresti), 14 medalii de argint: *Titiu Radu* (Universitatea din Bucuresti), Barzu Mihai (Departamentul de Informatică, Universitatea "Alexandru Ioan Cuza", Iași), Burtea Cosmin (Universitatea "Alexandru Ioan Cuza", Iași), Filip Laurian (Universitatea din București), Fodor Dan (Universitatea "Alexandru Ioan Cuza", Iași), Hlihor Petru (Universitatea din București), Sârbu Paul (Universitatea Tehnică "Gheorghe Asachi", Iași), Popescu Roxana-Irina (Universitatea din București), Mesaros Ionuț (Universitatea Tehnică, Cluj-Napoca), *Tucă Laurențiu* (Universitatea Politehnica București), *Gîlcă* Dragoş (Departamentul de Informatică, Universitatea "Alexandru Ioan Cuza", Iași), Beltic Marius Jimy Emanuel (Universitatea "Alexandru Ioan Cuza", Iași), Sasu Robert (Universitatea Politehnica București), Pleșca Iulia (Universitatea "Alexandru Ioan Cuza", Iași) și 19 medalii de bronz: Cervicescu Virgil (Academia Tehnică Militară, București), Raicea Marina (Academia Tehnică Militară, București), Mihăilă Ștefan (Departamentul de Informatică, Universitatea "Alexandru Ioan Cuza", Iași), Munteanu Alexandra-Irina (Universitatea din București), Vasile Mihaela Andreea (Universitatea Politehnica București), Mocanu Maria-Cristina (Departamentul de Informatică, Universitatea "Alexandru Ioan Cuza", Iași), Petre Luca (Universitatea Politehnica Bucuresti), Rublea Alina (Universitatea Politehnica București), Bobeș Maria Alexandra (Universitatea "Babeș Bolyai", Cluj-Napoca), Genes Cristian (Universitatea Tehnică "Gheorghe Asachi", Iași), Kolumban Jozsef (Universitatea "Babeş-Bolyai", Cluj-Napoca), Craus Sabina (Universitatea "Alexandru Ioan Cuza", Iași), Damanian Lavinia-Mariana (Departamentul de Informatică, Universitatea "Alexandru Ioan Cuza", Iași), Moldovan Dorin Vasile (Universitatea Tehnică, Cluj-Napoca), Mincu Diana (Universitatea Politehnica București), Birghila Corina (Universitatea "Ovidius", Constanța), Vlad Ilinca (Universitatea Tehnică "Gh. Asachi", Iași), Sav Adrian-Gabriel (Universitatea "Alexandru Ioan Cuza", Iași), Alban Andrei (Universitatea din București).

Echipa Universității Politehnica din București a fost formată din 8 studenți, dintre care 6 au reușit să obțină medalii.

Proba de concurs a constat în rezolvarea a 4 probleme pe durata a 5 ore. Vom prezenta mai jos soluțiile problemelor, precum și comentariile de rigoare.

2. Probleme, Soluții și Comentarii

Problema 1. Pentru un întreg dat $n \ge 1$, fie $f : [0,1] \to \mathbb{R}$ o funcție crescătoare. Demonstrați că

$$\int_{0}^{1} f(x) \mathrm{d}x \le (n+1) \int_{0}^{1} x^{n} f(x) \mathrm{d}x.$$

Găsiți toate funcțiile continue crescătoare pentru care egalitatea are loc.

* * *

Soluția 1. Pentru n dat și $x, y \in [0, 1]$ integrăm pe intervalul [0, 1] în raport cu x și în raport cu y inegalitatea evidentă

$$(x^{n} - y^{n})(f(x) - f(y)) \ge 0,$$
(1)

şi obţinem

$$\int_{0}^{1} \left(\int_{0}^{1} (x^{n} - y^{n})(f(x) - f(y)) \mathrm{d}x \right) \mathrm{d}y \ge 0,$$

sau

$$\int_{0}^{1} x^{n} f(x) dx - \int_{0}^{1} y^{n} dy \int_{0}^{1} f(x) dx - \int_{0}^{1} x^{n} dx \int_{0}^{1} f(y) dy + \int_{0}^{1} y^{n} f(y) dy \ge 0,$$

care ne dă inegalitatea cerută.

Dacă f este continuă, când egalitatea din enunț are loc, în (1) are loc de asemenea egalitate, deci funcția f trebuie să fie constantă.

Soluția a 2-a. Prin schimbare de variabilă avem

$$(n+1)\int_{0}^{1} x^{n} f(x) dx = \int_{0}^{1} f(\sqrt[n+1]{t}) dt$$

Cum f este crescătoare, rezultă că $f(x) \leq f(\sqrt[n+1]{x}), x \in [0,1]$, de unde prin integrare obținem inegalitatea cerută.

Soluția a 3-a. Problema este cazul particular al inegalității lui Cebâșev: Fie $f_1, f_2, \ldots, f_n : [a, b] \to \mathbb{R}$ funcții integrabile, pozitive, monotone.

1) Dacă f_1, f_2, \ldots, f_n sunt fie toate monoton crecătoare, fie toate monoton descrescătoare, atunci

$$\int_{a}^{b} f_1(x) \mathrm{d}x \cdot \int_{a}^{b} f_2(x) \mathrm{d}x \cdot \int_{a}^{b} f_n(x) \mathrm{d}x \le (b-a)^{n-1} \int_{a}^{b} f_1(x) f_2(x) \dots f_n(x) \mathrm{d}x.$$

2) Dacă f_1, f_2, \ldots, f_n sunt de monotonii diferite, atunci inegalitatea este de semn contrar.

În cazul nostru, pentru n = 2 și $f_1(x) = f(x), f_2(x) = x^n$ obținem concluzia din enunț.

Soluția a 4-a. Notăm
$$C = \int_{0}^{1} f(x) dx$$
 și fie $g : [0,1] \to \mathbb{R}$,
 $g(x) = f(x) - C$. Deci $\int_{0}^{1} g(x) dx = 0$.

Astfel, inegalitatea noastră este echivalentă cu $0 \leq \int_{0}^{1} x^{n}g(x) \mathrm{d}x$. Cum g
este crescătoare, rezultă că există $d = \sup\{x \in [0,1] : g(x) \leq 0\}$. Vom obține

$$\int_{0}^{1} x^{n} g(x) dx = -\int_{0}^{a} x^{n} |g(x)| dx + \int_{d}^{1} x^{n} |g(x)| dx$$

 $\operatorname{Cum} \int_{0}^{1} g(x) \mathrm{d}x = 0, \text{ avem} \int_{0}^{d} |g(x)| \mathrm{d}x = \int_{d}^{1} |g(x)| \mathrm{d}x. \text{ Înmulțind cu } d^{n},$

rezultă următorul șir de inegalități evidente

$$\int_{0}^{d} x^{n} |g(x)| \mathrm{d}x \le \int_{0}^{d} d^{n} |g(x)| \mathrm{d}x = \int_{d}^{1} d^{n} |g(x)| \mathrm{d}x \le \int_{d}^{1} x^{n} |g(x)| \mathrm{d}x,$$

de unde concluzia problemei noastre.

Comentariu. Problema a fost considerată de mulți concurenți drept una foarte ușoară. Cu toate acestea, au existat destui care n-au reușit să obțină punctajul maxim. De asemenea trebuie remarcat și faptul că au existat concurenți care nu și-au mai amintit sau pur și simplu nu știau inegalitatea lui *Cebâșev*.

Problema 2. Fie $A = (a_{ij})$ o matrice reală cu n linii și n coloane astfel încât $A^n \neq 0$ și $a_{ij}a_{ji} \leq 0$ pentru orice i, j. Demonstrați că există două numere nereale printre valorile proprii ale lui A.

Ivan Feshchenko, Kiev, Ucraina Soluție. Fie $\lambda_1, \ldots, \lambda_n$ toate valorile proprii ale acestei matrice. Polinomul caracteristic asociat matricei A este

$$P(\lambda) = \det(\lambda I_n - A) = \lambda^n - a_1 \lambda^{n-1} + \ldots + (-1)^n a_n$$

unde $a_1 = \sum_{i=1}^{} a_{ii}, a_n = \det A$ și a_i suma minorilor principali de ordiniai lui

A; de exemplu, $a_2 = \sum_{i,j} \begin{vmatrix} 0 & a_{ij} \\ a_{ji} & 0 \end{vmatrix}$.

Din ipoteză rezultă că $a_{ii} = 0$, deci $\sum_{k=1}^{n} \lambda_k = 0$. Suma $\sum_{i < j} \lambda_i \lambda_j = a_2$, i. e. $\sum_{i < j} \lambda_i \lambda_j = -\sum_{i < j} a_{ij} a_{ji}$. Aşadar, suma $\sum_{k=1}^{n} \lambda_k^2 = 2 \sum_{i < j} a_{ij} a_{ji} \le 0$.

Cum $A^n \neq 0$, A are cel puțin o valoarea proprie nenulă, de unde rezultă că cel puțin o valoare proprie nu este reală. Deoarece polinomul caracteristic al lui A are coeficienți reali și o valoare proprie este complexă, rezultă că și

conjugata complexă a acestei valori este valoare proprie pentru A, deci ceea ce trebuia demonstrat.

Comentariu. Problema a fost considerată de dificultate medie. Mai mult de jumătate dintre concurenți n-au reușit să o rezolve complet. Din păcate mulți dintre studenții români au întâmpinat dificultăți mari în abordarea acestei probleme. Acest lucru arată că este nevoie de mai multă Algebră Liniară în programele noastre din universități.

Problema 3. Fie vectorii $\overline{a}, \overline{b}, \overline{c} \in \mathbb{R}^n$. Arătați că

$$(||\overline{a}||\langle \overline{b}, \overline{c} \rangle)^2 + (||\overline{b}||\langle \overline{a}, \overline{c} \rangle)^2 \le ||\overline{a}||||\overline{b}||(||\overline{a}||||\overline{b}|| + |\langle \overline{a}, \overline{b} \rangle|)||\overline{c}||^2,$$

unde $\langle \overline{x}, \overline{y} \rangle$ este produsul scalar al vectorilor $\overline{x}, \overline{y}$ şi $\|\overline{x}\|^2 = \langle \overline{x}, \overline{x} \rangle$.

Dan Schwarz, București, România

Soluția 1. (Soluția trigonometrică.) Vom demonstra mai întâi că, dați vectorii $u, v \in \mathbb{R}^n$, cu ||u|| = ||v|| = 1 (deci $u, v \in S^{n-1}$), atunci

$$\sup_{|x||=1} \left(\langle u, x \rangle^2 + \langle v, x \rangle^2 \right) = 1 + |\langle u, v \rangle|.$$

Considerăm reprezentarea vectorială unicăx=x'+x'' cu $x'\in {\rm span}(u,v)$ și $x''\perp {\rm span}(u,v).$ Atunci

$$1 = ||x||^2 = \langle x' + x'', x' + x'' \rangle =$$

 $= ||x'||^2 + 2\left\langle x', x'' \right\rangle + ||x''||^2 = ||x'||^2 + ||x''||^2 \ (\text{relația lui } Pitagora),$

deci $||x'|| \le 1$. Avem $|\langle u, x \rangle| = |\langle u, x' \rangle| \le |\langle u, y \rangle|$, unde y = 0 dacă x' = 0 şi $y = \frac{x'}{||x'||}$ dacă $x' \ne 0$ (deci ||y|| = 1). Analog $|\langle v, x \rangle| = |\langle v, x' \rangle| \le |\langle v, y \rangle|$. Maximul este aşadar obținut când $x \in \text{span}(u, v)$.

Deci problema noastră vectorială a fost transferată în spațiul de dimensiune 2, cu vectorii unitari u, v și x. Capetele vectorilor $\pm u$ și $\pm v$ împart cercul S^1 în patru arce, fiecare măsurând cel mult π (și posibil 0, când $u = \pm v$); capătul lui x va fi într-unul din ele. Fie ω măsura acelui arc și α, β , cu $\alpha + \beta = \omega$, măsurile arcelor dintre capătul lui x și capetele acelui arc. Atunci $\langle u, x \rangle^2 + \langle v, x \rangle^2 = \cos^2 \alpha + \cos^2 \beta$, din binecunoscuta interpretare geometrică a produsului scalar (independent de dimensiunea spațiului). Atunci

$$\langle u, x \rangle^2 + \langle v, x \rangle^2 = \cos^2 \alpha + \cos^2 \beta = 1 + \frac{1}{2} (\cos 2\alpha + \cos 2\beta) =$$
$$= 1 + \cos(\alpha + \beta) \cos(\alpha - \beta) \le 1 + |\cos \omega| = 1 + |\langle u, v \rangle|,$$

cu egalitate în cazurile evidente:

• când $\alpha = \beta = \omega/2$, deci când $|\langle u, x \rangle| = |\langle v, x \rangle|$, aşadar

$$x = (\pm u \pm v)/|| \pm u \pm v||,$$

unde semnele sunt astfel încât $0 \le \omega < \frac{\pi}{2}$;

• când $\omega = \pi/2$, deci când $\langle u, v \rangle = 0$, i.e. $u \perp v$, pentru orice $x \in \operatorname{span}(u, v)$.

Atunci, pentru orice \overline{a} , \overline{b} , \overline{c} nenuli, luăm $u = \frac{\overline{a}}{||\overline{a}||}$, $v = \frac{\overline{b}}{||\overline{b}||}$, $x = \frac{\overline{c}}{||\overline{c}||}$, aşadar

$$\langle u, x \rangle^2 = \frac{1}{||\overline{a}|| \cdot ||\overline{c}||^2} \langle \overline{a}, \overline{c} \rangle^2, \langle v, x \rangle^2 = \frac{1}{||\overline{b}|| \cdot ||\overline{c}||^2} \langle \overline{b}, \overline{c} \rangle^2, \langle u, v \rangle = \frac{\langle \overline{a}, \overline{b} \rangle}{||\overline{a}|| \cdot ||\overline{b}||},$$

deci relația demonstrată devine

$$\frac{1}{||\overline{a}||} \langle \overline{a}, \overline{c} \rangle^2 + \frac{1}{||\overline{b}||} \left\langle \overline{b}, \overline{c} \right\rangle^2 \leq \left(1 + \frac{|\langle \overline{a}, \overline{b} \rangle|}{||\overline{a}|| \cdot ||\overline{b}||} \right) ||\overline{c}||^2,$$

echivalent cu inegalitatea cerută, de asemenea adevărată pentru $\overline{a},\overline{b}$ sau \overline{c} egale cu zero.

Soluția a 2-a. (Soluție cu forme pătratice) Pentru||x|| = ||u|| = ||v|| = 1 avem

$$0 \le ||\lambda x + \mu u + \nu v||^2 = \langle \lambda x + \mu u + \nu v, \lambda x + \mu u + \nu v \rangle = \lambda^2 + \mu^2 + \nu^2 + 2\lambda \mu \langle x, u \rangle + 2\lambda \nu \langle x, v \rangle + 2\mu \nu \langle u, v \rangle,$$

o formă pătratică care ia valori pozitive pentru orice parametri reali λ, μ, ν , deci corespunzând unei matrice pozitiv semidefinită

$$\left[\begin{array}{ccc} 1 & \langle x, u \rangle & \langle x, v \rangle \\ \langle x, u \rangle & 1 & \langle u, v \rangle \\ \langle x, v \rangle & \langle u, v \rangle & 1 \end{array}\right].$$

Minorii principali de ordinul 1 sunt așadar pozitivi, ceea ce ne dă semipozitivitatea normei; minorii principali de ordinul 2 sunt pozitivi, ceea ce ne dă inegalitatea *Cauchy-Buniakowski-Schwarz*, $1 \ge \langle u, v \rangle^2$; de asemenea, determinatul matricii este pozitiv

$$\Delta = 1 - \left(\langle u, v \rangle^2 + \langle x, u \rangle^2 + \langle x, v \rangle^2 \right) + 2 \langle u, v \rangle \langle x, u \rangle \langle x, v \rangle \ge 0,$$

care poate fi scris

$$\langle x, u \rangle^2 + \langle x, v \rangle^2 \le 1 + |\langle u, v \rangle| - |\langle u, v \rangle|(1 + |\langle u, v \rangle| - 2|\langle x, u \rangle \langle x, v \rangle|).$$

Dar $\langle x, u \rangle^2 + \langle x, v \rangle^2 \ge 2 |\langle x, u \rangle \langle x, v \rangle|$, care înlocuit în relație ne dă

$$(1 - |\langle u, v \rangle|)(1 + |\langle u, v \rangle| - 2|\langle x, u \rangle\langle x, v \rangle|) \ge 0.$$

Apoi, fie 1 = $|\langle u, v \rangle|$, când 1 + $|\langle u, v \rangle|$ = 2 \geq 2 $|\langle x, u \rangle \langle x, v \rangle|$ (din inegalitatea *Cauchy-Buniakowski-Schwarz*), fie 1 > $|\langle u, v \rangle|$, ceea ce implică 1 + $|\langle u, v \rangle| \geq$ 2 $|\langle x, u \rangle \langle x, v \rangle|$. Aşadar, întot
deauna

$$\langle x, u \rangle^2 + \langle x, v \rangle^2 \le 1 + |\langle u, v \rangle|.$$

Soluția a 3-a (Soluția cu multiplicatorii lui Lagrange). Definim

$$L(x,\lambda) = \langle u, x \rangle^{2} + \langle v, x \rangle^{2} - \lambda(||x||^{2} - 1)$$

și considerăm sistemul

$$\frac{\partial L}{\partial x_i} = 2u_i \langle u, x \rangle + 2v_i \langle v, x \rangle - 2\lambda x_i = 0, \text{ pentru } 1 \le i \le n,$$

şi

$$\frac{\partial L}{\partial \lambda} = ||x||^2 - 1 = 0.$$

Avem

$$0 = \frac{1}{2} \sum_{i=1}^{n} x_i \frac{\partial L}{\partial x_i} = \langle u, x \rangle \sum_{i=1}^{n} u_i x_i + \langle v, x \rangle \sum_{i=1}^{n} v_i x_i - \lambda \sum_{i=1}^{n} x_i^2 = \langle u, x \rangle^2 + \langle v, x \rangle^2 - \lambda ||x||^2 = \langle u, x \rangle^2 + \langle v, x \rangle^2 - \lambda.$$

Pe de altă parte,

$$0 = \frac{1}{2} \sum_{i=1}^{n} u_i \frac{\partial L}{\partial x_i} = \langle u, x \rangle \sum_{i=1}^{n} u_i^2 + \langle v, x \rangle \sum_{i=1}^{n} v_i u_i - \lambda \sum_{i=1}^{n} x_i u_i = 0$$

$$= \langle u, x \rangle ||u||^2 + \langle v, x \rangle \langle u, v \rangle - \lambda \langle u, x \rangle = \langle u, x \rangle + \langle u, v \rangle \langle v, x \rangle - \lambda \langle u, x \rangle,$$

și analog pentruv, deci obținem sistemul de două ecuații (în variabilele $\langle u,x\rangle$ și $\langle v,x\rangle)$

$$\left\{ \begin{array}{l} \left(1-\lambda\right)\left\langle u,x\right\rangle +\left\langle u,v\right\rangle\left\langle v,x\right\rangle =0,\\ \left\langle u,v\right\rangle\left\langle u,x\right\rangle +\left(1-\lambda\right)\left\langle v,x\right\rangle =0. \end{array} \right. \right.$$

Determinantul Δ al matricei acestui sistem este $\Delta = (1 - \lambda)^2 - \langle u, v \rangle^2$, și dacă este nenul, unica soluție este cea trivială $\langle u, x \rangle = \langle v, x \rangle = 0$, când expresia noastră atinge un minim evident egal cu zero (atunci $\lambda = 0$, aceasta înseamnă $\langle u, v \rangle \neq \pm 1$, care s-ar traduce în $u \neq \pm v$ și $x \perp \text{span}(u, v)$). Aşadar, suntem interesați de $\Delta = 0$ (pentru toate celelalte puncte critice), ducând la $\lambda = 1 \pm \langle u, v \rangle$, deci $\lambda = 1 + |\langle u, v \rangle|$ într-un punct de maxim și $\lambda = 1 - |\langle u, v \rangle|$ într-un punct critic. Există câteva cazuri particulare.

Când $u = \pm v$, deci $\langle u, v \rangle = \pm 1$, situația este simplă. Avem maximul 2 când $x = \pm u$ și minimul 0 când $x \perp u$.

Când $u \perp v$, deci $\langle u, v \rangle = 0$, atunci λ (deci și expresia) este 1 în punctele de maxim, pentru orice $x \in \text{span}(u, v)$, și λ (deci și expresia) este 0 în punctele de minim, pentru orice $x \perp \text{span}(u, v)$.

Soluția a 4-a. Vom reduce problema la cazul \mathbb{R}^3 . Considerăm o bază ortonormală pentru care primii trei vectori generează subspațiul vectorial generat de $\{\overline{a}, \overline{b}, \overline{c}\}$. O transformare ortogonală păstrează norma și produsul scalar, deci este suficient să demonstrăm inegalitatea în noua bază în care \overline{a} , \overline{b} , \overline{c} au primele trei componente nenule.

În \mathbb{R}^3 vom folosi un argument de geometrie. Mai întâi, observăm că orice schimbare de semn a vectorilor \overline{a} , \overline{b} , \overline{c} nu schimbă semnul inegalității inițiale.

Fie $\alpha = \sphericalangle(\overline{b}, \overline{c}), \ \beta = \sphericalangle(\overline{c}, \overline{a})$ şi $\gamma = \sphericalangle(\overline{a}, \overline{b})$. Aranjăm semnele astfel încât $\alpha, \beta \in \left[0, \frac{\pi}{2}\right]$ şi considerăm triedrul *OABC* având laturile *OA*, *OB*, OC determinate de vectorii $\overline{a}, \overline{b}, \overline{c}$ astfel încât proiectând vectorul \overline{c} pe planul OAB, proiecția sa să rămână în interiorul unghiului OAB.

Este cunoscut faptul că, dacă x, y, z satisfac x + y = z, atunci $\cos^2 x + \cos^2 y + \cos^2 z - 2\cos x \cos y \cos z = 1$. Astfel, avem de arătat că $\cos^2 \alpha + \cos^2 \beta \le 1 + |\cos \gamma|$. Se observă că, dacă unul dintre unghiurile α sau β este egal cu $\frac{\pi}{2}$, atunci inegalitatea este evidentă.

Proiectăm punctul C pe planul OAB în punctul P și considerăm $\alpha_1 =$ $\triangleleft POB, \beta_1 = \triangleleft POA.$ Proiectăm P pe OA, respectiv OB în punctele A_1 , respectiv B_1 . În triunghiurile dreptunghice OB_1C și OB_1P avem

$$\cos^2 \alpha = \frac{OB_1^2}{OC^2} \le \frac{OB_1^2}{OP^2} = \cos^2 \alpha_1.$$

Analog avem $\cos^2 \beta \le \cos^2 \beta_1$. Din $\alpha_1 + \beta_1 = \gamma$ obţinem

$$\cos^2 \alpha_1 + \cos^2 \beta_1 = 1 + 2 \cos \alpha_1 \cos \beta_1 \cos \gamma - \cos^2 \gamma.$$

Arătăm că $1 + 2\cos\alpha_1\cos\beta_1\cos\gamma - \cos^2\gamma \le 1 + \cos\gamma$.

Dacă $\cos \gamma = 0$, atunci inegalitatea este evidentă. Dacă $\cos \gamma \neq 0$, inegalitatea este echivalentă cu

 $2\cos\alpha_1\cos\beta_1 \le 1 + \cos\gamma = \cos(\alpha_1 + \beta_1) = 1 + \cos\alpha_1\cos\beta_1 - \sin\alpha_1\sin\beta_1,$ care în final va da $\cos(\alpha_1 - \beta_1) \leq 1$.

Soluția a 5-a (Nicolae Beli). Asemănător cu soluția 4, presupunem $\alpha,\beta\in\Big[0,\frac{\pi}{2}\Big].$ Trebuie arătat că

$$|\cos\gamma| \ge \cos^2\alpha + \cos^2\beta - 1 = \frac{1}{2}(\cos 2\alpha + \cos 2\beta) = \cos(\alpha + \beta)\cos(\alpha - \beta).$$

Cum $|\alpha - \beta| \leq \frac{\pi}{2}$, aven $\cos(\alpha - \beta) \geq 0$. Prin urmare, putem presupune că $\cos(\alpha - \beta), \cos(\alpha + \beta) > 0$. Insă α, β, γ sunt unghiurile triedrului OABC(care poate fi și degenerat), deci avem $0 \le \gamma \le \alpha + \beta \le \pi$. Rezultă că $0 < \cos(\alpha + \beta) \le \cos \gamma$, care împreună cu $0 < \cos(\alpha - \beta) \le 1$ implică $\cos(\alpha + \beta)\cos(\alpha - \beta) \le \cos\gamma.$

Comentariu. Această problemă s-a dovedit a fi cea mai grea problemă din concurs, doar câțiva reușind să o rezolve complet. Dintre români, doar Victor Pădureanu a rezolvat problema. În rezolvarea problemei de față s-au folosit aceleași ingrediente ca într-o altă problemă, un pic mai veche, dată la ultimul test de selecție a lotului olimpic al României din anul 2007. Ea suna cam așa:

Problemă înrudită. Pentru $n \ge 2$ întreg pozitiv, se consideră numerele reale $a_i, b_i, i = \overline{1, n}$, astfel încât

$$\sum_{i=1}^{n} a_i^2 = \sum_{i=1}^{n} b_i^2 = 1, \sum_{i=1}^{n} a_i b_i = 0.$$

Să se arate că

$$\left(\sum_{i=1}^{n} a_i\right)^2 + \left(\sum_{i=1}^{n} b_i\right)^2 \le n.$$

Pentru soluții și comentarii recomandăm cititorilor Romanian Mathematical Competitions 2007, paginile 97–100.

Problema 4. Fie $f : [0,1] \to \mathbb{R}$ o funcție crescătoare de clasă C^2 . Definim șirurile date de $L_n = \frac{1}{n} \sum_{k=0}^{n-1} f\left(\frac{k}{n}\right)$ și $U_n = \frac{1}{n} \sum_{k=1}^n f\left(\frac{k}{n}\right)$, $n \ge 1$. Intervalul $[L_n, U_n]$ se împarte în trei segmente egale. Demonstrați că, pentru n suficient de mare, numărul $I = \int_0^1 f(x) dx$ aparține segmentului din mijloc dintre aceste trei segmente egale.

Alexander Kukush, Kiev, Ucraina Soluția 1. Enunțăm și demonstrăm mai întâi următoarea lemă : Lema. Pentru $f \in C^2[0,1]$: $L_n = I - \frac{f(1) - f(0)}{2n} + O\left(\frac{1}{n^2}\right), \quad n \to \infty.$ Demonstrație. Notăm C = f(1) - f(0). Considerăm

$$I - L_n = \sum_{k=0}^{n-1} \int_{\frac{k}{n}}^{\frac{k+1}{n}} \left(f(x) - f\left(\frac{k}{n}\right) \right) \mathrm{d}x =$$

$$=\sum_{k=0}^{n-1}\int_{\frac{k}{n}}^{\frac{k+1}{n}} \left(f'\left(\frac{k}{n}\right)\left(x-\frac{k}{n}\right) + \frac{1}{2}f''(\theta_{kn}^x)\left(x-\frac{k}{n}\right)^2\right) dx =$$
$$=\frac{1}{2n^2}\sum_{k=0}^{n-1}f'\left(\frac{k}{n}\right) + r_n,$$
(1)

unde

$$|r_n| \le \frac{1}{2} \max|f''| \cdot \sum_{k=0}^{n-1} \frac{1}{3} \left(x - \frac{k}{n}\right)^3 \Big|_{\frac{k}{n}}^{\frac{k+1}{n}} \le \frac{\text{const}}{n^2}.$$

Aici $\theta_{kn}^x \in \left[\frac{k}{n}, \frac{k+1}{n}\right]$ sunt punctele intermediare din teorema lui *Taylor* în vecinătatea punctului $\frac{k}{n}$.

Analog

$$\int_{0}^{1} f' \mathrm{d}x - \frac{1}{n} \sum_{k=0}^{n-1} f'\left(\frac{k}{n}\right) = \sum_{k=0}^{n-1} \int_{\frac{k}{n}}^{\frac{k+1}{n}} f''(\widetilde{\theta}_{kn}^{x})\left(x - \frac{k}{n}\right) \mathrm{d}x = O\left(\frac{1}{n}\right).$$

Deci în partea dreaptă a relației (1) avem

$$\frac{1}{n}\sum_{k=0}^{n-1}f'\left(\frac{k}{n}\right)\to\int_{0}^{1}f'\mathrm{d}x=C,\quad n\to\infty,$$

și eroarea în partea dreaptă a relației (1) este de ordin $O\left(\frac{1}{n^2}\right)$.

Deci din (1) avem

_

$$I - L_n = \frac{C}{2n} + O\left(\frac{1}{n^2}\right), \quad n \to \infty.$$

Acum ne întoarcem la soluția problemei noastre și avem $U_n = L_n + \frac{C}{n}$. Fie $x_n = L_n + \frac{kC}{3n}$, k = 1, 2. Atunci

$$x_n = I + \frac{C}{n} \left(\frac{k}{3} - \frac{1}{2}\right) + O\left(\frac{1}{n^2}\right).$$

Pentru k = 1 avem $\frac{k}{3} - \frac{1}{2} < 0$, deci $x_n < I$ pentru n suficient de mare; pentru k = 2 avem $\frac{k}{3} - \frac{1}{2} > 0$, deci $x_n > I$ pentru n suficient de mare. Aşadar, pentru n suficient de mare

$$L_n + \frac{C}{3n} < I < L_n + \frac{2C}{3N}.$$

Soluția a 2-a. Fie $F(t) = \int_{0}^{t} f(x) dt$. Atunci avem

$$I = \int_{0}^{1} f(x) dx = F(1) - F(0) =$$

F(1) - F($\frac{n-1}{n}$) + F($\frac{n-1}{n}$) - F($\frac{n-2}{n}$) + ... + F($\frac{1}{n}$) - F(0)

Din formula lui *Taylor* de ordinul 2, vom obține (pentru punctele 0, $\frac{1}{n}, \frac{2}{n}, \dots, \frac{n-1}{n}$): $F(x) = F\left(\frac{k}{n}\right) + \left(x - \frac{k}{n}\right)F'\left(\frac{k}{n}\right) + \frac{1}{2}\left(x - \frac{k}{n}\right)^2F''(\theta),$ cu $\theta \in \left[\frac{k}{n}, x\right], k = \overline{0, n-1},$ sau $F(x) = F\left(\frac{k}{n}\right) + \left(x - \frac{k}{n}\right)f\left(\frac{k}{n}\right) + \frac{1}{2}\left(x - \frac{k}{n}\right)^2f'(\theta).$ Pentru $x = \frac{1}{n}, \dots, 1$ avem $F\left(\frac{k+1}{n}\right) - F\left(\frac{k}{n}\right) = \frac{1}{n}f\left(\frac{k}{n}\right) + \frac{1}{2n^2}f'(\theta_k),$ cu $\theta \in \left[\frac{k}{n}, \frac{k+1}{n}\right], k = \overline{0, n-1}$

cu $\theta_k \in \left[\frac{k}{n}, \frac{k+1}{n}\right], \ k = \overline{0, n-1}.$

Însumând aceste relații pentru k = 0, 1, n - 1 obținem

$$I = L_n + \frac{1}{2n}\sigma_n,$$

unde σ_n este suma *Riemann* pentru $\int_0^1 f'(x)dx = f(1) - f(0).$
Intervalul din mijloc este $\left[\frac{2}{3}L_n + \frac{1}{3}U_n, \frac{1}{3}L_n + \frac{2}{3}U_n\right] = [u_n, v_n].$
Dacă $f(0) = f(1)$ atunci f este constantă, ceea ce nu se poate

Dacă f(0) = f(1), atunci \tilde{f} este constantă, ceea ce nu se poate. Presupunem că f(1) > f(0).

Avem $n(I-u_n) = \frac{n}{3}(L_n - U_n) + \frac{1}{2}\sigma_n = \frac{1}{3}(f(0) - f(1)) + \frac{1}{2}\sigma_n$ şi trecând la limită când $n \to \infty$, obținem $n(I - u_n) \to \frac{1}{6}(f(1) - f(0)) > 0$, de unde pentru n suficient de mare, $I > u_n$.

În același mod, avem

$$n(v_n - I) = \frac{2n}{3}(U_n - L_n) - \frac{1}{2}\sigma_n \to \frac{1}{6}(f(1) - f(0)) > 0.$$

Remarcă. Relația $I = L_n + \frac{1}{2n}\sigma_n$ mai poate fi dedusă, din următorul fapt:

Fie $f:[0,1] \to \mathbb{R}$ derivabilă cu derivata integrabilă pe [0,1]. Atunci

$$\lim_{n \to \infty} n \left[\int_{0}^{1} f(x) dx - \frac{1}{n} \sum_{k=1}^{n} f\left(\frac{k}{n}\right) \right] = \frac{f(0) - f(1)}{2}.$$

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 $Demonstrație. \ {\rm Fie} \ \Delta = \left(0,\frac{1}{n},\frac{2}{n},\ldots,\frac{n-1}{n},1\right) \ {\rm o} \ {\rm diviziune} \ {\rm a} \ {\rm intervalului} \ [0,1]. \ {\rm Not\ \mbox{am}}$

$$M_k = \sup_{x \in \left[\frac{k-1}{n}, \frac{k}{n}\right]} f'(x), \qquad m_k = \inf_{x \in \left[\frac{k-1}{n}, \frac{k}{n}\right]} f'(x), \quad k = \overline{1, n}$$

Atunci avem

$$E_n = n \left[I - \frac{1}{n} \sum_{k=1}^n f\left(\frac{k}{n}\right) \right] = n \sum_{k=1}^n \left[\int_{\frac{k-1}{n}}^{\frac{k}{n}} f(x) dx - \frac{1}{n} f\left(\frac{k}{n}\right) \right] =$$
$$= n \sum_{k=1}^n \left[\int_{\frac{k-1}{n}}^{\frac{k}{n}} \left(f(x) - f\left(\frac{k}{n}\right) \right) dx \right].$$

Aplicând teorema lui *Lagrange* pe intervalul $\left[x, \frac{k}{n}\right]$, rezultă că există $c_k(x) \in \left(\frac{k-1}{n}, \frac{k}{n}\right)$ astfel încât $f(x) - f\left(\frac{k}{n}\right) = \left(x - \frac{k}{n}\right) f'(c_k(x))$. Atunci $E_n = n \sum_{k=1}^n \int_{\frac{k-1}{n}}^{\frac{k}{n}} \left(x - \frac{k}{n}\right) f'(c_k(x)) dx$. Dar $m_k \leq f'(c_k(x)) \leq M_k$, deci avem $-\frac{1}{2n} \sum_{k=1}^n M_k \leq E_n \leq -\frac{1}{2n} \sum_{k=1}^n m_k$. Thereford by limitation physical limitation $E_n = \frac{f(0) - f(1)}{2n}$

Trecând la limită obținem $\lim_{n \to \infty} E_n = \frac{f(0) - f(1)}{2}$.

Comentariu. Problema a fost considerată una foarte dificilă, iar puțini concurenți au reușit să o rezolve în întregime, cu toate că problema putea fi abordată și de un elev foarte bine pregătit de clasa a XII-a. Dintre concurenții români, doar câțiva au reușit să obțină mai mult de jumătate din punctaj.

3. Concluzii

Olimpiada pentru studenți SEEMOUS din acest an a fost o reușită pentru studenții români, care s-au descurcat onorabil. Mulți dintre concurenți au arătat destul de bine pregătiți, dovada fiind medaliile obținute. Pe de altă parte trebuie remarcat și faptul că parcurgerea unor noțiuni din anul I de facultate în liceu poate constitui un real avantaj pentru cei care erau foarte buni la matematică în liceu.

În ceea ce privește studenții Universității Politehnica din București, putem spune că au avut o comportare decentă. Din păcate niciunul dintre ei n-a reușit să obțină medalie de aur, cu toate că *Laurențiu Țucă* a fost foarte aproape dacă nu ar fi existat abordarea superficială a ultimei probleme. De asemenea, trebuie tras un semnal de alarmă în ceea ce privește rezultatele la prima problemă. Doar doi din cei 8 studenți români și-au amintit de inegalitatea lui *Cebâșev* (forma discretă sau integrală), deși ea fusese făcută la una din pregătiri într-o formă sau alta. La prima problemă, o altă soluție corectă a fost dată de *Alina Rublea*, care a obținut aproape punctajul maxim având mici scăpări în tratarea cazului de egalitate.

Problema 2 a fost cea mai abordată de studenții Politehnicii. Laurențiu Țucă a fost cel mai aproape de o soluție completă. Din păcate, el nu a văzut un argument simplu care l-ar fi ajutat să finalizeze. Ceilalți au abordat problema, însă au existat mici deficiențe în înțelegerea noțiunilor de polinom caracteristic și valoare proprie. Din păcate, trebuie să semnalăm din nou faptul că din cauza comprimării materiei în anul I studenții, fie ei chiar și cei mai buni, nu reușesc să-și însușească noțiunile importante din anumite capitole ale algebrei liniare.

Problema 3 a fost cea mai dificilă din concurs și a fost rezolvată complet doar de 3 concurenți, printre care și *Victor Pădureanu*, fost component al lotului olimpic și al celui lărgit pe perioada liceului. O altă încercare, dar nefinalizată, a avut *Robert Sasu*. La fel, parcă se pune prea puțin accent la cursul de Algebră pe aceste noțiuni de vectori, respectiv forme pătratice, determinanți *Gram*, inegalitățile *Cauchy-Schwarz*, *Bessel*, etc.

Problema 4 a fost a doua din concurs ca dificultate. Unii dintre studenții români s-au descurcat bine, iar dintre cei ai Politehnicii doar *Laurențiu Țucă* a fost aproape de a lua punctaj maxim, dacă ar fi demonstrat lema de mai sus.

In concluzie, suntem de părere că, din cauza cantității mari de materie care a fost comprimată aproape într-un singur semestru, studenții nu mai au posiblitatea să mai simtă importanța și mai ales gustul unor noțiuni peste care se trece cu mult prea mare ușurință la orele de curs și seminar. Graba predării materiei nu va avea ca rezultat decât familiarizarea insuficientă chiar și a celor mai buni studenți cu noțiunile noi de la cursurile Analiză Matematică, Algebră Liniară și Superioară, Geometrie. Din păcate, aceasta este o meteahnă care datează de mult prea mulți ani în învățământul matematic din România.

PROBLEMS

Authors should submit proposed problems to office@rms.unibuc.ro or to gmaproblems@gmail.com. Files shold be in PDF or DVI format. Once a problem is accepted and considered for publication, the author will be asked to submit the TeX file also. The referee process will usually take between several weeks and two months. Solutions may also be submitted to the same e-mail adress. For this issue, solutions should arrive before November 15, 2011.

Editors: Mihai Cipu, Radu Gologan, Călin Popescu, Dan Radu Assistant Editor: Cezar Lupu

PROPOSED PROBLEMS

323. Let \mathcal{C} be the set of the circles in the plane and \mathcal{L} be the set of the lines in the plane. Show that there exist bijective maps $f, g : \mathcal{C} \to \mathcal{L}$ such that for any circle $C \in \mathcal{C}$, the line f(C) is tangent at C and the line g(C) contains the center of C.

Proposed by Marius Cavachi, Ovidius University of Constanța, Romania.

324. Consider the set

$$K:=\left\{\,f\left(\sqrt[4]{20}\,,\sqrt[6]{500}\,\right)\mid f(X,Y)\in\mathbb{Q}[X,Y]\,\right\}.$$

- (a) Show that K is a field with respect to the usual addition and multiplication of real numbers.
- (b) Find all the subfields of K.
- (c) If one considers K as a vector space $\mathbb{Q}K$ over the field \mathbb{Q} in the usual way, find the dimension of $\mathbb{Q}K$.
- (d) Exhibit a vector space basis of $_{\mathbb{O}}K$.

Proposed by Toma Albu, Simion Stoilow Institute of Mathematics of the Romanian Academy, Bucharest, Romania.

325. We call *toroidal chess board* a regular chess board (of arbitrary dimension) in which the opposite sides are identified in the same direction. Show that the maximum number of kings on a toroidal chess board of dimensions $m \times n$ $(m, n \in \mathbb{N})$ such that each king attacks no more than six other kings is less than or equal to $\frac{4mn}{5}$ and the inequality is sharp.

Proposed by Eugen Ionașcu, Columbus State University, Columbus, GA, USA.

326. For
$$t > 0$$
 define $H(t) = \sum_{n=0}^{\infty} \frac{t^n}{n!(n+1)!}$. Show that
$$\lim_{t \to \infty} \frac{t^{3/4}H(t)}{\exp(2\sqrt{t})} = \frac{1}{2\sqrt{\pi}}.$$

Proposed by Moubinool Omarjee, Jean Lurçat High School, Paris, France.

327. Let $f : [a, b] \to \mathbb{R}$ be a convex and continuous function. Prove that:

a)
$$\mathcal{M}(a;b) + (b-a)f\left(\frac{a+b}{2}\right) \ge \mathcal{M}\left(\frac{3a+b}{4};\frac{3b+a}{4}\right);$$

b) $3\mathcal{M}\left(\frac{2a+b}{3};\frac{2b+a}{3}\right) + \mathcal{M}(a;b) \ge 4\mathcal{M}\left(\frac{3a+b}{4};\frac{3b+a}{4}\right).$
Here $\mathcal{M}(x,y) = \frac{1}{y-x}\int_{-\infty}^{y} f(t)dt.$

Proposed by Cezar Lupu, Polytechnic University of Bucharest, Bucharest, Romania, and Tudorel Lupu, Decebal High School, Constanța, Romania.

328. Given any positive integers m, n, prove that the set $\{1, 2, 3, \ldots, m^{n+1}\}$ can be partitioned into m subsets A_1, A_2, \ldots, A_m , each of size m^n , such that

$$\sum_{a_1 \in A_1} a_1^k = \sum_{a_2 \in A_2} a_2^k = \dots = \sum_{a_m \in A_m} a_m^k, \text{ for all } k = 1, 2, \dots, n.$$

Proposed by Cosmin Pohoață, student Princeton University, Princeton, NJ, USA.

329. Let $p \ge 11$ be prime number. Show that, if

$$\sum_{j=1}^{(p-1)/2} \frac{1}{j^6} = \frac{a}{b}$$

with a, b relatively prime, then p divides a.

Proposed by Marian Tetiva, Gheorghe Roşca Codreanu National College, Bârlad, Romania.

330. Determine all nonconstant monic polynomials $f \in \mathbb{Z}[X]$ such that $\varphi(f(p)) = f(p-1)$ for all natural prime numbers p. Here $\varphi(n)$ is the *Euler* totient function.

Proposed by Vlad Matei, student Cambridge University, Cambridge, UK.

331. Let $\mathcal{B}_n = \{(x_1, \ldots, x_n) \in \mathbb{R}^n \mid x_i \leq x_{i+2} \text{ for } 1 \leq i \leq n-2\}$ and let $\mathcal{B} = \bigcup_{n \geq 1} \mathcal{B}_n$. On \mathcal{B} we define the relation \leq as follows. If $x, y \in \mathcal{B}$, $x = (x_1, \ldots, x_m)$ and $y = (y_1, \ldots, y_n)$, we say that $x \leq y$ if $m \geq n$ and for any $1 \leq i \leq n$ we have either $x_i \leq y_i$ or 1 < i < m and $x_i + x_{i+1} \leq y_{i-1} + y_i$. Prove that (\mathcal{B}, \leq) is a partially ordered set.

Proposed by Nicolae Constantin Beli, Simion Stoilow Institute of Mathematics of the Romanian Academy, Bucharest, Romania.

332. For a positive integer $n = \prod_{i=1}^{s} p_i^{\alpha_i}$ denote by $\Omega(n) := \sum_{i=1}^{s} \alpha_i$ the

total number of prime factors of n (counting multiplicities). Of course, by default $\Omega(1) = 0$. Define now $\lambda(n) := (-1)^{\Omega(n)}$, and consider the sequence $\mathfrak{S} := (\lambda(n))_{n \geq 1}$. You are asked to prove the following claims on \mathfrak{S} :

a) It contains infinitely many terms $\lambda(n) = -\lambda(n+1)$;

b) It is not ultimately periodic;

c) It is not ultimately constant over an arithmetic progression;

- d) It contains infinitely many pairs $\lambda(n) = \lambda(n+1)$;
- d) It contains infinitely many terms $\lambda(n) = \lambda(n+1) = 1$;
- e) It contains infinitely many terms $\lambda(n) = \lambda(n+1) = -1$.

Proposed by Dan Schwarz, Bucharest, Romania.

333. Show that there do not exist polynomials $P, Q \in \mathbb{R}[X]$ such that

$$\int_{0}^{\log \log n} \frac{P(x)}{Q(x)} dx = \frac{1}{p_1} + \frac{1}{p_2} + \ldots + \frac{1}{p_n}, n \ge 1,$$

where p_n is the *n*th prime number.

Proposed by Cezar Lupu, Polytechnic University of Bucharest, Bucharest and Cristinel Mortici, Valahia University of Târgovişte, Târgovişte, Romania.

334. Let a, b be two positive integers with a even and $b \equiv 3 \pmod{4}$. Show that $a^m + b^m$ does not divide $a^n - b^n$ for any even $m, n \geq 3$.

Proposed by Octavian Ganea, student École Polytechnique Fédérale de Lausanne, Lausanne, Switzerland.

335. Let m and n be positive integers with $m \leq n$ and let $A \in \mathcal{M}_{m,n}(\mathbb{R})$ and $B \in \mathcal{M}_{n,m}(\mathbb{R})$ be matrices such that rank $A = \operatorname{rank} B = m$. Show that there exists $C \in \mathcal{M}_n(\mathbb{R})$ such that $A \cdot C \cdot B = I_m$, where I_m denotes the m by m unit matrix.

Proposed by Vasile Pop, Technical University Cluj-Napoca, Cluj-Napoca, Romania.

336. Show that the sequence $(a_n)_{n\geq 1}$ defined by $a_n = [2^n\sqrt{2}] + [3^n\sqrt{3}]$, $n \geq 1$, contains infinitely many odd numbers and infinitely many even numbers. Here [x] is the integer part of x.

Proposed by Marius Cavachi, Ovidius University of Constanţa, Constanţa, Romania.

SOLUTIONS

295. Determine all nonconstant polynomials $P \in \mathbb{Z}[X]$ such that P(p) is square-free for all prime numbers p.

Proposed by Vlad Matei, student University of Bucharest, Bucharest, Romania.

Solution by the author. We will need the following observation. Lemma. For all nonconstant polynomials $f \in \mathbb{Z}[X]$ the set

$$A = \{q \text{ prime } | \exists p \text{ prime, } q | f(p) \}$$

is infinite.

Proof of the lemma. Assume the contrary. Let $A = \{q_1, q_2, \ldots, q_n\}$ and let $M = q_1q_2\cdots q_n$. Let k be an arbitrary natural number. According to Dirichlet's theorem, the arithmetic progression M^kr+1 with $r \in \mathbb{N}^*$ contains an infinity of prime numbers.

Let us note that $f(M^k r + 1) \equiv f(1) \pmod{q_i^k}$, $\forall i = 1, \ldots, n$. Below we will assume that the leading coefficient of f is positive, the other case is similar considering -f. We know that $\lim_{x\to\infty} (f(x) - x) = \infty$, thus for rsufficiently large $f(M^k r + 1) \geq M^k r$. Since A is finite, it follows that there is an index j with $q_j^k | f(M^k r + 1)$. We would have that $q_j^k | f(1)$. Now for ksufficiently large $q_i^k > |f(1)|, \forall i = 1, \ldots, n$, and we are done unless f(1) = 0. In this case the lemma is obvious, since f(p) = (p-1)g(p), with $g \in \mathbb{Z}[X]$.

First of all, we claim that f(0) = 0. Assume the contrary $f(0) \neq 0$. We can choose infinitely many primes q such that $f(q) = d \cdot p_1 p_2 \cdots p_m$ with $p_i \neq p_j$ for $1 \leq i \neq j \leq m$ and (f(0), f(q)) = d. We know from *Taylor* expansion of polynomials that

$$f(q+tp_i) = f(q) + f'(q) \cdot tp_i + f''(q) \cdot \frac{(tp_i)^2}{2!} + \dots \equiv f(q) + f'(q) \cdot tp_i \pmod{p_i^2}$$

for $1 \leq i \leq m$.

If we would have $p_n \nmid f'(q)$ for a certain index n, then we could choose t such that $f(q) + f'(q) \cdot tp_n \equiv 0 \pmod{p_i^2}$, since this is equivalent to $f'(q) \cdot t \equiv -\frac{f(q)}{p_n} \pmod{p_i^2}$ and f'(q) is invertible modulo p_n . This means that there is $m \in \mathbb{Z}$ such that $p_n^2 \mid f(m)$. Now $m = q + tp_n$, and if $(m, p_n) > 1$ it follows that $p_n \mid q$. Since $f(q) = d \cdot p_1 p_2 \dots p_m$, we have that $p_n \mid f(q)$, so $p_n \mid (q, f(q))$. But $(q, f(q)) \mid f(0)$, and we deduce that $p_n \mid d$, a contradiction. Thus $(m, p_n) = 1$, and by Dirichlet's theorem we could find a prime r in the arithmetic progression $m + ap_n^2$, $a \in \mathbb{N}^*$. We would get that $f(r) \equiv f(m) \equiv 0 \pmod{p_i^2}$, in contradiction with the hypothesis.

Therefore there is no such index and it follows $p_1 p_2 \cdots p_m |f'(q)$. This means $|f'(q)| \ge |p_1 p_2 \cdots p_m|$, which is equivalent to $\frac{|f'(q)|}{|f(q)|} \ge \frac{1}{d} \ge \frac{1}{f(0)}$. To finish the proof, we make q arbitrarily large and it results that

To finish the proof, we make q arbitrarily large and it results that $\lim_{x\to\infty} \frac{f'(x)}{f(x)} \ge \frac{1}{f(0)}$. This contradicts the well known fact that $\lim_{x\to\infty} \frac{f'(x)}{f(x)} = 0$. This contradiction proves that our initial claim is true, so that f(0) = 0. If we write $f(X) = X^i g(X), g(0) \neq 0$, we get immediately from the hypothesis that i = 1 and g is constant. So f(X) = cX. Morever, if c has a prime factor l, then f(l) is not square-free. Thus $c = \pm 1$.

We conclude that the only such polynomials are f(X) = -X and f(X) = X.

296. Let $a_1, a_2, \ldots, a_n, b_1, b_2, \ldots, b_n > 0$ and $x_1, x_2, \ldots, x_n, y_1, y_2, \ldots, y_n$ be real numbers. Show that

$$\left(\sum_{1\leq i,j\leq n} x_i x_j \min(a_i, a_j)\right) \left(\sum_{1\leq i,j\leq n} y_i y_j \min(b_i, b_j)\right) \geq \sum_{1\leq i,j\leq n} x_i y_j \min(a_i, b_j).$$

Proposed by Alin Gălățan, student University of Bucharest, Bucharest, Romania.

Solution by the editors. Let λ_A be the characteristic function of an arbitrary set A. Consider the functions $f, g: [0, \infty) \to \mathbb{R}$ defined by

$$f(x) = \sum_{i=1}^{n} x_i \lambda_{[0,a_i]}(x), \quad g(x) = \sum_{i=1}^{n} y_i \lambda_{[0,b_i]}(x)$$

We have

$$\int_{0}^{\infty} f^{2}(x) \mathrm{d}x = \sum_{1 \le i,j \le n} x_{i} x_{j} \int_{0}^{\infty} \lambda_{[0,a_{i}]}(x) \lambda_{[0,b_{i}]}(x) \mathrm{d}x = \sum_{1 \le i,j \le n} x_{i} x_{j} \min(a_{i},a_{j}).$$

Analogously, we obtain

$$\int_{0}^{\infty} g^2(x) \mathrm{d}x = \sum_{1 \le i,j \le n} y_i y_j \min(b_i, b_j), \quad \int_{0}^{\infty} f(x) g(x) \mathrm{d}x = \sum_{1 \le i,j \le n} x_i y_j \min(a_i, b_j).$$

Finally, our inequality reduces to

$$\int_{0}^{\infty} f^{2}(x) \mathrm{d}x \int_{0}^{\infty} g^{2}(x) \mathrm{d}x \ge \left(\int_{0}^{\infty} f(x)g(x) \mathrm{d}x\right)^{2},$$

which is nothing else than the celebrated integral version of the Cauchy-Schwarz inequality. $\hfill \Box$

297. Let S^2 be the bidimensional sphere and $\alpha > 0$. Show that for any positive integer n and $A_1, A_2, \ldots, A_n, B_1, B_2, \ldots, B_n$ and C_1, C_2, \ldots, C_n arbitrary points on S^2 , there exists $P_n \in S^2$ such that

$$\sum_{i=1}^{n} P_n A_i^{\alpha} = \sum_{i=1}^{n} P_n B_i^{\alpha} = \sum_{i=1}^{n} P_n C_i^{\alpha}$$

if and only if $\alpha = 2$.

Proposed by Marius Cavachi, Ovidius University of Constanța, Constanța, Romania.

Solution by the author. We will divide the proof in two cases $\alpha = 2$ and $\alpha \neq 2$. In the first case let us define the function $f: S^2 \to \mathbb{R}^2$,

$$f(P) = \left(\sum_{i=1}^{n} PA_i^2 - \sum_{i=1}^{n} PB_i^2, \sum_{i=1}^{n} PA_i^2 - \sum_{i=1}^{n} PC_i^2\right)$$

Since f(-P) = -f(P), $\forall P \in S^2$, according to *Borsuk-Ulam*'s Theorem, there is a $P_n \in S^2$ such that $f(P_n) = (0,0)$.

We can provide an elementary argument for this fact, as it follows.

Let H_1, H_2 be the geometrical locus of the points P in space for which $\sum_{i=1}^{n} PA_i^2 = \sum_{i=1}^{n} PB_i^2, \text{ respectively } \sum_{i=1}^{n} PA_i^2 = \sum_{i=1}^{n} PC_i^2. H_1 \text{ and } H_2 \text{ are planes}$

which both pass trough the center of the sphere. In consequence, $H_1 \cap H_2$ contains a line, which passes through the center of the sphere, and we can choose P_n as one of its intersections with the surface of the sphere.

For the second case, when $\alpha \neq 2$, we will use standard cartesian coordinates, and let N = (0, 0, 1) and S = (0, 0, -1), the standard north and south pole of the sphere. We denote with S^1 the standard plane section given by the ecuation z = 0.

For *n* arbitrary and $i \in \{1, 2, ..., n\}$ we choose $A_i = N$, $B_i = S$, and C_i are the vertices of a regular polygon with *n* sides circumscribed around S_1 . A point P_n which satisfies the statement of the problem must be on S^1 and satisfies $\sum_{i=1}^{n} P_n C_i^{\alpha} = n(\sqrt{2})^{\alpha}$.

Also, since $\lim_{n\to\infty} \frac{1}{n} \sum_{i=1}^{n} P_n C_i^{\alpha} = \frac{1}{2\pi} \int_{P \in S^1} Q P^{\alpha} ds = I$, where Q is for

example the point (1, 0, 0), we deduce that

$$I = (\sqrt{2})^{\alpha}.\tag{1}$$

Now

$$I = \frac{1}{2\pi} \int_{0}^{2\pi} \left(2\sin\left(\frac{t}{2}\right) \right)^{\alpha} dt = \frac{2^{\alpha}}{\pi} \int_{0}^{\pi} \sin^{\alpha} u \, du =$$

$$= \frac{2^{\alpha}}{\pi} \int_{0}^{\frac{\pi}{2}} (\sin^{\alpha} v + \cos^{\alpha} v) \, \mathrm{d}v = \frac{2^{\alpha}}{2} (\sin^{\alpha} v_{0} + \cos^{\alpha} v_{0}),$$

using the mean theorem, and $v_0 \in \left(0; \frac{\pi}{2}\right)$.

We deduce that

$$\frac{I}{2^{\frac{\alpha}{2}}} = \frac{1}{2} \left[(2\sin^2(v_0))^{\frac{\alpha}{2}} + (2\cos^2(v_0))^{\frac{\alpha}{2}} \right] = \frac{1}{2} [(1+w)^{\beta} + (1-w)^{\beta}],$$

unde $\beta = \frac{\alpha}{2}, w \in (0; 1)$. If $\alpha > 2$, then $\beta > 1$, and using the Bernoulli inequality we have $I = \frac{1}{2} - \frac{1}{2}$. $(1+w)^{\beta} + (1-w)^{\beta} > 1 + \beta w + 1 - \beta w > 2$. Then $\frac{I}{2^{\frac{\alpha}{2}}} > \frac{1}{2} \cdot 2 = 1$, a contradiction with (1).

Similarly for
$$\alpha \in (0; 2)$$
, we obtain the contradiction $\frac{1}{2^{\frac{\alpha}{2}}} < 1$.

Remark. Marius Olteanu had a geometrical approach based on Leibniz relations.

298. Let t be an odd number. Find all monic polynomials $P \in \mathbb{Z}[X]$ such that for all integers n there exists an integer m for which

$$P(m) + P(n) = t.$$

Proposed by Octavian Ganea, student École Polytechnique Fédérale de Lausanne, Lausanne, Switzerland and Cristian Tălău, student Polytechnic University of Bucharest, Bucharest, Romania.

Solution by the authors. Firstly we prove that the degree of P is odd. If ${\cal P}$ would have even degree, then it would be bounded from below, but from our statement taking $n \to \infty$ we have that $\lim_{n \to +\infty} P(n) = +\infty$, it follows that there is a sequence of values of m for which $P(m) \to -\infty$, which contradicts the fact that P is bounded from below.

So we have proven that degree of P is odd. Now let M be large enough such that P is strictly increasing on $[M, +\infty)$ (it is possible since P' has a finite number of roots and $\lim_{x \to +\infty} P'(x) = +\infty$). So from our hypothesis, it follows that P is strictly increasing on $(-\infty, N]$ (*), where N is such P(M) + P(N) = t.

Let $m \ge M$ and $n \in \mathbb{Z}$ such that P(m) + P(-n) = t. From the hypothesis there is a $k \in \mathbb{Z}$ such that P(m+1) + P(-n-k) = t. It is clear from (*) that $k \ge 1$. Also from the hypothesis there is a $r \ge 1$ with the property that P(m+r) + P(-n-1) = t. If $r \ge 2$ or $k \ge 2$, then P(m+1) + P(-n-k) < P(m+r) + P(-n-1), a contradiction. Thus P(m+1) + P(-n-1) = t, so by induction P(m+k) + P(-n-k) = t, $\forall k \in \mathbb{N}.$

Let a = m - n and $Q(X) = P\left(\frac{a}{2} + X\right) + P\left(\frac{a}{2} - X\right) - t$. Let us notice that it has infinitely many roots of the form $x = \frac{m+n}{2} + k$, $\forall k \in \mathbb{N}$. It follows that Q = 0. For x = 0 we have $P\left(\frac{a}{2}\right) = \frac{t}{2} \notin \mathbb{Z}$, so a is odd. Let us notice that if we denote $k = \deg(P)$ we have $2^{k-1}P\left(\frac{a}{2}\right) = 2^{k-2}t$ and for $k \ge 3$ we know that $2^{k-2}t \in \mathbb{Z}$, so $2^{k-1}P\left(\frac{a}{2}\right) \in \mathbb{Z}$ and if we write the expression of P, it would follow that $\frac{a^k}{2} \in \mathbb{Z}$, which is in contradiction with a odd. Thus k = 1 so P(X) = X + b, which verifies the hypothesis. \Box

299. Find all functions $f : \{1, 2, ...\} \rightarrow \{1, 2, ...\}$ satisfying the following properties:

i) a - b divides f(a) - f(b) for all $a, b \in \{1, 2, \ldots\}$;

ii) if a, b are relatively prime, so are f(a) and f(b).

Proposed by Gabriel Dospinescu, École Normale Supérieure de Paris, Paris, France and Fedor Nazarov, University of Wisconsin, Madison, WI, USA.

Solution by the authors. First, we claim that any prime factor of f(n) divides n. Assume that p divides f(n) and does not divide n. Since p divides f(n+p) - f(n), p also divides f(n+p), so it divides gcd(f(n), f(n+p)). As gcd(n, n+p) = 1, we also have gcd(f(n), f(n+p)) = 1, which makes the last result impossible. The first claim is thus proved.

Write $f(p) = p^{g(p)}$ for each prime p, for some $g(p) \ge 0$. We will prove that g is constant. Fix odd primes p, q and a positive integer m and define $u = \frac{p^{2^m} + 1}{2}$ and $v = \frac{q^{2^n} + 1}{2}$ for some n such that gcd(u, v) = 1 (this is possible since classically any prime factor of v is at least 2^n , so if n is large enough, v will be relatively prime to u). The Chinese Remainder Theorem combined with *Dirichlet*'s theorem give us a prime r such that $r \equiv p \pmod{u}$, $r \equiv q \pmod{v}$. Therefore, u divides $f(r) - f(p) = r^{g(r)} - p^{g(p)}$. Since u also divides $r^{g(r)} - p^{g(r)}$, it follows that u divides $p^{|g(r)-g(p)|} - 1$. Recalling the definition of u, it immediately follows that 2^{m+1} divides g(r) - g(p), and doing the same with v yields that 2^{m+1} also divides g(q) - g(r), so that 2^{m+1} divides g(p) - g(q). Since m was arbitrary, we must have that g is constant, say g(p) = d for any prime p.

Finally, if n is a positive integer we have $n - p|f(n) - f(p) = f(n) - p^d$ and $n - p|n^d - p^d$, so that $n - p|f(n) - n^d$. Since this holds for any prime p, we must have $f(n) = n^d$ for all n. Obviously, all these functions are solutions to the problem. **300.** Consider the sequence $(a_n)_{n\geq 1}$ defined by $a_1 = 2$ and $a_{n+1} = 2a_n + \sqrt{3(a_n^2 - 1)}$, $n \geq 1$. Show that the terms of a_n are positive integers and any odd prime p divides $a_p - 2$.

Proposed by Alin Gălățan, student University of Bucharest and Cezar Lupu, student University of Bucharest, Bucharest, Romania.

Solution by the authors. It is well-known that the function $ch : \mathbb{R} \to (1,\infty)$, defined by $ch(x) = \frac{e^x + e^{-x}}{2}$, is onto and thus there exists $\alpha \ge 1$ such that $ch(\alpha) = 2$, which implies that $sh(\alpha) = \sqrt{3}$, with $sh : \mathbb{R} \to [1,\infty)$ defined by $sh(x) = \frac{e^x - e^{-x}}{2}$. It follows that

$$a_{n+1} = \operatorname{ch}(\alpha)a_n + \operatorname{sh}(\alpha)\sqrt{a_n^2 - 1}, \quad \forall n \ge 1.$$

We shall prove by induction that $a_n = ch(n\alpha)$, $\forall n \ge 1$. The case n = 1is obvious. We assume that $a_n = ch(n\alpha)$ for some $n \ge 1$ and we prove that $a_{n+1} = ch((n+1)\alpha)$. Indeed, since $ch(n\alpha + \alpha) = ch\alpha \cdot ch(n\alpha) + sh\alpha \cdot sh(n\alpha)$, $\forall n \ge 1$, we obtain that $a_n = ch(n\alpha)$, $\forall n \ge 1$.

In what follows, we prove by induction that all terms a_n are positive integers. Indeed, we assume that for some $n \ge 2$ one has $a_k \in \mathbb{N}$ for all integers $1 \le k < n$, and we prove that $a_n \in \mathbb{N}$, too. Assume by contradiction that $a_n \notin \mathbb{N}^*$. This means that there exists a square-free number d > 1 such that $a_n = M + N\sqrt{d}$, so that $a_n \notin \mathbb{Q}$. We have

$$2^{2n} = \operatorname{ch}^{n}(\alpha) = \left(\frac{e^{\alpha} + e^{-\alpha}}{2}\right)^{n},$$

which is equivalent to

$$2^{2n} = \sum_{k=0}^{n} \binom{n}{k} e^{(n-k)\alpha - k\alpha} = \sum_{k=0}^{n} \binom{n}{k} \left(e^{(n-2k)\alpha} + e^{-(n-2k)\alpha} \right) = \sum_{k=0}^{n} \binom{n}{k} a_{|n-2k|}.$$

Since a_k is integer for all k < n, it follows that $a_n \in \mathbb{Q}$, which contradicts the initial assumption. Thus we get $a_n \in \mathbb{N}$ for all $n \ge 1$.

Finally, consider a prime number p > 2. We have

$$4^{p} = \sum_{k=0}^{p} \binom{k}{p} \cdot e^{(p-2k)\alpha} = \sum_{k=0}^{\frac{p-1}{2}} \binom{p}{k} \left(e^{(p-2k)\alpha} + e^{-(p-2k)\alpha} \right) = \sum_{k=0}^{\frac{p-1}{2}} \binom{p}{k} \cdot 2a_{p-2k}.$$

Since p divides $\binom{p}{k}$ and, from *Fermat*'s Little Theorem, $4^p \equiv 4 \pmod{p}$, it follows that $2 \equiv a_p \pmod{p}$ and hence the conclusion follows immediately having in view that p is odd.

Solution by Marian Tetiva, Gheorghe Roşca Codreanu National College, Bârlad, Romania. That a_n is positive for every $n \ge 1$ immediately follows by induction, and this readily implies $a_{n+1} > a_n$ for all n (thus, the sequence is strictly increasing).

By squaring $a_{n+1}-2a_n = \sqrt{3(a_n^2-1)}$ one gets $a_{n+1}^2-4a_{n+1}a_n+a_n^2 = -3$ for all positive integers n; therefore,

$$a_{n+1}^2 - 4a_{n+1}a_n + a_n^2 = a_n^2 - 4a_na_{n-1} + a_{n-1}^2, \ \forall n \ge 2.$$

This one can be read as

$$(a_{n+1} - a_{n-1})(a_{n+1} + a_{n-1} - 4a_n) = 0,$$

hence the terms of $(a_n)_{n\geq 1}$ also verify the recurrence

$$a_{n+1} - 4a_n + a_{n-1} = 0, \ \forall n \ge 2$$

(as $a_{n+1} - a_{n-1}$ is always nonzero – actually positive). A canonical induction shows now that a_n is an integer for each $n \ge 1$ (as it starts with $a_1 = 2$, $a_2 = 7$, and $a_{n+1} = 4a_n - a_{n-1}$, for $n \ge 2$).

Also canonical is to find the formula of the general term; we have

$$a_n = \frac{1}{2} \left(\left(2 + \sqrt{3} \right)^n + \left(2 - \sqrt{3} \right)^n \right), \ \forall n \ge 1.$$

Now we have the following general result: if x_1, x_2, \ldots, x_m are the zeros of a monic polynomial with integer coefficients, and p is a positive prime number, then

$$(x_1 + x_2 + \dots + x_m)^p - (x_1^p + x_2^p + \dots + x_m^p)$$

is an integer which is divisible by p.

Indeed, the first part follows easily via the fundamental theorem of symmetric polynomials. By the multinomial formula and the fact that each multinomial coefficient $\frac{p!}{i_1!i_2!\cdots i_m!}$ is divisible by p whenever none of i_1, i_2, \ldots, i_m equals p,

$$\frac{1}{p}\left((x_1 + x_2 + \dots + x_m)^p - (x_1^p + x_2^p + \dots + x_m^p)\right)$$

is a sum of products of powers of x_1, x_2, \ldots, x_m with integer coefficients, therefore it is an algebraic integer. Being rational and algebraic integer, the above number is an integer, which was to be proven.

(See also *Example 6* at page 191 of *Problems from the Book* by T. Andreescu and G. Dospinescu, XYZ Press, 2008. And note that yet another beautiful solution, which uses matrices, is given to that problem starting with page 192.)

In our case one can consider $x_1 = 2 + \sqrt{3}$, and $x_2 = 2 - \sqrt{3}$; the above theorem implies that

$$4^p - \left(\left(2+\sqrt{3}\right)^p + \left(2-\sqrt{3}\right)^p\right)$$

is divisible by p for every positive prime number p. If p is odd, from here we infer that

$$x_p - 2^{2p-1} = \frac{1}{2} \left(\left(2 + \sqrt{3} \right)^p + \left(2 - \sqrt{3} \right)^p \right) - 2^{2p-1}$$

is divisible by p. Then

$$x_p - 2 = (x_p - 2^{2p-1}) + (2^{2p-1} - 2)$$

is also divisible by p, because $2^{2p-1} - 2$ is divisible by p, too, according to *Fermat*'s Little Theorem (which is, somehow, generalized by the above mentioned result).

Note that the fact that $x_p - 2$ is divisible by p can also be obtained by using the binomial development in the formula for x_p .

Remark. Similar solutions with the second approach were also given by *Nicuşor Minculete, Braşov* and *Marius Olteanu, Râmnicu Vâlcea.*

301. Determine the greatest prime number p = p(n) such that there exists a matrix $X \in SL(n, \mathbb{Z})$ with $X^p = I_n$ and $X \neq I_n$.

Proposed by Victor Vuletescu, University of Bucharest, Bucharest, Romania.

Solution by Marian Tetiva, Gheorghe Roşca Codreanu National College, Bârlad, Romania. We show that p(n) is the greatest prime which does not exceed n + 1.

First we prove that $p(n) \leq n+1$. Suppose p is a prime, and that there exists $X \in SL(n,\mathbb{Z})$ with $X^p = I_n$. Thus X is a root of the polynomial

$$T^{p} - 1 = (T - 1)(T^{p-1} + \dots + T + 1)$$

and then the minimal polynomial of X over \mathbb{Q} is one of T-1, $T^{p-1}+\cdots+T+1$ (both factors are irreducible), or T^p-1 itself. Since we look for $X \neq I_n$, the first case is excluded. Because the degree of the minimal polynomial of a matrix of order n is at most n, we find that either $p-1 \leq n$ or $p \leq n$; in both cases $p \leq n+1$ follows.

Now we show that for p = p(n), the greatest prime which does not exceed n + 1, there is a matrix $X \in SL_n(\mathbb{Z})$ with $X^p = I_n$. As $p - 1 \leq n$, we can consider $X = (x_{ij})_{1 \leq i,j \leq n}$ with $x_{11} = x_{12} = \cdots = x_{1,p-1} = -1$, $x_{21} = x_{32} = \cdots = x_{p-1,p-2} = 1$, and all other entries equal to 0. For this matrix we have $X^{p-1} + \cdots + X + I_n = 0_n$ (it is built starting from the companion matrix of $T^{p-1} + \cdots + T + 1$), therefore $X^p - I_n = 0_n \Leftrightarrow X^p = I_n$, too (just multiply the previous equality with $X - I_n$).

Note that when $p(n) \leq n$ (that is, when n + 1 is not a prime), a permutation matrix corresponding to a cycle of order p can be used as yet an example.

302. Let $n \geq 2$ and denote by $\mathcal{D} \subset M_n(\mathbb{C})$ and $\mathcal{C} \subset M_n(\mathbb{C})$ the set of diagonal and circulant matrices, respectively. Consider $\mathcal{V} = [D, C] =$

= span{ $dc - cd \mid d \in \mathcal{D}, c \in \mathcal{C}$ }. Prove that \mathcal{C}, \mathcal{D} , and \mathcal{V} span $M_n(\mathbb{C})$ if and only if n is prime.

Proposed by Remus Nicoară, University of Tennessee, Knoxville, TN, USA.

Solution by the author. We begin with a standard notation; $E_{i,j}$ is the matrix with 0 or 1 as entries and which has the value 1 only on the position (i, j). Let us note firstly that the set \mathcal{D} is generated by the matrices $E_{h,h} = D_h$, where $1 \le h \le n$. The set \mathcal{C} is generated by the matrices X^l ,

with
$$l \le n$$
, where $X = \begin{pmatrix} 0 & 1 & 0 & \dots & \dots & 0 \\ 0 & 0 & 1 & \dots & \dots & 0 \\ \dots & \dots & \dots & \dots & \dots & \dots \\ 0 & 0 & 0 & \dots & 1 & 0 \\ 1 & 0 & 0 & \dots & \dots & 0 \end{pmatrix}$, so $X^{l} = \sum_{i=1}^{n} E_{i,i+l}$, where

the indices are taken modulo n.

For $1 \leq k \leq n$ and $l \leq n$, we have that $[D_k, X^l] = E_{k,k+l} - E_{k-l,k}$ generate the space \mathcal{V} . It is easy to verify that $\mathcal{C}, \mathcal{D} \perp \mathcal{V}$, where the scalar product is given by $\langle X, Y \rangle = \operatorname{Tr}(\overline{Y^t}X)$.

Since span(\mathcal{C}, \mathcal{D}) has dimension 2n - 1, it follows that span($\mathcal{C}, \mathcal{D}, \mathcal{V}$) = $= M_n(\mathbb{C})$ if and only if dim_{$\mathbb{C}} V = (n-1)^2$.</sub>

Let $V_{h,l} = E_{h,h+l} - E_{h-l,h}$. They span the space \mathcal{V} . We must look at the dependence relations between $V_{h,l}$. Let $c_{h,l} \in \mathbb{C}$ such that $\sum_{h,l} c_{h,l}V_{h,l} = O_n$. This is equivalent to $\sum_{i,j} (c_{i,j-i} - c_{j-i,j})E_{i,j} = O_n$, so $c_{i,j-i} = c_{j-i,j}$.

If we denote d = j - i, this means $c_{i,d} = c_{i+d,d}$, for any i, d, so we deduce $c_{i,d} = c_{i+md,d}$, for any $m \ge 1$. Now if $d \ne 0$ and n is prime the set $\{i + md\}$, where $m \ge 1$, contains all the residues modulo n, besides i, and since $i \ne j$ we can find an m_0 such that $i + m_0 d \equiv j \pmod{p}$. Thus $c_{i,d} = c_{j,d}, \forall i, j, d \neq 0$.

Now let us look at the linear transformation $c_{h,l} \rightarrow \sum_{h,l} c_{h,l} V_{h,l}$. From

what we obtained, for n prime, we know that its kernel is determined, on the first n-1 columns, by the values on the first row, and the last column is arbitrary. Thus we have 2n-1 "degrees of freedom", so the kernel has dimension 2n-1. This means that the image has dimension $n^2 - 2n + 1 =$ $=(n-1)^2.$

Finally, let us see what happens when n is not prime. We prove that the kernel has a bigger dimension, and the proof ends. This is easy to see since $\{i + md\}$ cannot span the whole residues when d is not coprime with n, and we can choose such a $d \neq 0$ and d < n. In this case on the d-th column we have at least one more "degree of freedom" so the dimension of the kernel is bigger than 2n-1.

303. Prove that

$$\sum_{n=0}^{\infty} \frac{1}{\Gamma(n+3/2)} = \frac{2e}{\sqrt{\pi}} \int_{0}^{1} e^{-x^{2}} dx$$

and

$$\sum_{n=0}^{\infty} \frac{\Gamma(-n+1/2) \Gamma(n+1/2)}{\Gamma(n+3/2)} = \frac{\sqrt{\pi}}{2e} \int_{0}^{1} e^{x^{2}} dx$$

where $\Gamma(x) = \int_{0}^{\infty} e^{-t} t^{x-1} dt$ is *Euler*'s Gamma function.

Proposed by Cezar Lupu, student University of Bucharest, Bucharest and Tudorel Lupu, Decebal High School, Constanţa, Romania.

Solution by the authors. On the first serie, we use Series expansion of Incomplete Gamma function, namely

$$\Gamma(a-1) - \Gamma(a-1,x) = x^a \cdot \left(\frac{1}{(a-1)x} - \frac{1}{a} + \frac{x}{2(a+1)} - \frac{x^2}{6(a+2)} + \frac{x^3}{24(a+3)} - \frac{x^4}{120(a+4)} + \frac{x^5}{720(a+5)} + O\left(x^6\right)\right).$$

We multiply the equation by $x^{(1-a)}$ to make things more polynomial-like and thus we have

$$x^{1-a} \big(\Gamma(a-1) - \Gamma(a-1,x) \big) = \frac{1}{a-1} - \frac{x}{a} + \frac{x^2}{2(a+1)} - \frac{x^3}{6(a+2)} + \frac{x^4}{24(a+3)} - \frac{x^5}{120(a+4)} + O\left(x^6\right).$$

The magic happens when multiplying by e^x , it will make a translation in the x terms, or in series form

$$\begin{split} \left(\sum_{j=0}^{\infty} \frac{x^j}{j!}\right) \left(\sum_{i=0}^{\infty} \frac{(-1)^i x^i}{i!(a+i-1)}\right) &= \sum_{k=0}^{\infty} \sum_{i=0}^k \frac{(-1)^i x^k}{i!(a+i-1)(k-i)!}.\\ \text{So, we obtain}\\ \mathrm{e}^x x^{1-a} \big(\Gamma(a-1) - \Gamma(a-1,x)\big) &= \frac{1}{a-1} + \frac{x}{(a-1)a} + \frac{x^2}{(a-1)a(a+1)} + \\ &+ \frac{x^3}{(a-1)a(a+1)(a+2)} + \frac{x^4}{(a-1)a(a+1)(a+2)(a+3)} + O\left(x^5\right). \end{split}$$

Finally eliminating the term
$$a - 1$$
, we have

$$(a-1)e^{x}x^{1-a}\left(\Gamma(a-1) - \Gamma(a-1,x)\right) = 1 + \frac{x}{a} + \frac{x^{2}}{a(a+1)} + \frac{x^{3}}{a(a+1)(a+2)} + \frac{x^{3}}{a(a+2)} + \frac{x^{3}}{a(a+2)} + \frac{x^{3}}{a(a+2)} + \frac{x^{3}}{a(a+2)}$$

$$+\frac{x^3}{a(a+1)(a+2)} + \frac{x^4}{a(a+1)(a+2)(a+3)} + O\left(x^5\right) = \sum_{i=0}^{\infty} \frac{x^i(a-1)!}{(i+a-1)!}.$$

Dividing by $\Gamma(a)$ and using Gamma function instead of factorial, we deduce

$$\sum_{i=0}^{\infty} \frac{x^i}{\Gamma(a+i)} = \frac{(-1+a)\mathrm{e}^x x^{1-a} \left(\Gamma(-1+a) - \Gamma(-1+a,x)\right)}{\Gamma(a)}$$

The left and right hand side are valid for any positive real number a.

Making x = 1, $a = \frac{3}{2}$ and using erf function, i.e., $\operatorname{erf}(x) = \frac{2}{\pi} \int_{0}^{x} e^{-t^{2}} dt$,

instead of the incomplete Gamma, we conclude the solution.

For the second series, we just have to put $x \to -x$ and the solution is almost the same as above.

Remark. A similar solution, but a little bit more laborious, was given by *Marius Olteanu, Râmnicu Vâlcea*.

304. Let $f, g : [0, 1] \to \mathbb{R}$ be two functions such that f is continuous and g is increasing and differentiable, with $g(0) \ge 0$. Prove that if for any $t \in [0, 1]$

$$\int_{t}^{1} f(x) \mathrm{d}x \ge \int_{t}^{1} g(x) \mathrm{d}x,$$

then

$$\int_{0}^{1} f^{2}(x) \mathrm{d}x \geq \int_{0}^{1} g^{2}(x) \mathrm{d}x.$$

Proposed by Andrei Ciupan, student Harvard University, Boston, MA, USA.

Solution by the author. Let F and G be fixed, but otherwise arbitrary, primitives of f and g, respectively. From the AM-GM inequality we have

$$\int_{0}^{1} f^{2}(x) dx + \int_{0}^{1} g^{2}(x) dx \ge 2 \cdot \int_{0}^{1} f(x) g(x) dx.$$
(1)

By integration by parts we obtain

$$\int_{0}^{1} f(x)g(x)dx = F(x) \cdot g(x) \Big|_{0}^{1} - \int_{0}^{1} F(x)g'(x)dx.$$
(2)

The hypothesis tells us that $F(1) - F(x) \ge G(1) - G(x)$ for any $x \in [0, 1]$, which can be rewritten as $-F(x) \ge G(1) - F(1) - G(x)$. Since g

is differentiable and increasing, we have $g'(x) \ge 0$ for any x in [0,1]. From these two facts we obtain

$$-\int_{0}^{1} F(x)g'(x)dx \ge \int_{0}^{1} (G(1) - F(1))g'(x)dx - \int_{0}^{1} G(x)g'(x)dx \Leftrightarrow$$
$$\Leftrightarrow -\int_{0}^{1} F(x)g'(x)dx \ge (G(1) - F(1))(g(1) - g(0)) - G(x) \cdot g(x)\Big|_{0}^{1} + \int_{0}^{1} g^{2}(x)dx.$$

By combining this last relation with (2), we obtain

$$\int_{0}^{1} f(x)g(x)dx \ge F(1)g(1) - F(0)g(0) + (G(1) - F(1))(g(1) - g(0))$$

$$-G(1)g(1) + G(0)g(0) + \int_{0}^{1} g^{2}(x) \mathrm{d}x.$$

By reducing and grouping, we finally obtain

$$\int_{0}^{1} f(x)g(x)dx \ge g(0)\left(F(1) - F(0) + G(0) - G(1)\right) + \int_{0}^{1} g^{2}(x)dx.$$

Therefore, by taking into account relation (1), we obtain

$$\int_{0}^{1} f^{2}(x) dx \ge \int_{0}^{1} g^{2}(x) dx + 2g(0) \left(\int_{0}^{1} f(x) dx - \int_{0}^{1} g(x) dx \right) \ge \int_{0}^{1} g^{2}(x) dx,$$

since $g(0) \ge 0$ and $\int_{0}^{1} f(x) dx \ge \int_{0}^{1} g(x) dx.$

Remark. A similar solution, but a little bit more complicated, was given by *Marius Olteanu*, *Râmnicu Vâlcea*.

305. Let K be a field and $f \in K[X]$ with $\deg(f) = n \ge 1$ having distinct roots x_1, x_2, \ldots, x_n . For $p \in \{1, 2, \ldots, n\}$ let S_1, S_2, \ldots, S_p be the symmetric fundamental polynomials in x_1, x_2, \ldots, x_p . Show that

$$[K(S_1, S_2, \dots, S_p) : K] \le \binom{n}{p}$$

Proposed by Marius Cavachi, Ovidius University of Constanța, Constanța, Romania.

Solution by the author. We consider firstly the following fields, $F = K(S_1, \ldots, S_p), \ L = K(x_1, \ldots, x_p), \ M = K(x_1, x_2, \ldots, x_n).$ We have the inclusions $K \subset F \subset L \subset M$. The extension $K \hookrightarrow M$ is Galois and let $G = \operatorname{Gal}(M/K)$ be the *Galois* group, and finally let us denote with G_1 the image of the group $\operatorname{Gal}(M/F)$ through the canonical inclusion $\operatorname{Gal}(M/F) \hookrightarrow G$.

The group Ga(M/T) through the set of subfields of M and $\left|\frac{G}{\operatorname{Stab}(L)}\right| = |\operatorname{Orb}(L)|$. Since for $\sigma \in G, \sigma(L)$ is uniquely determined by the subset $\{\sigma(x_1), \ldots, \sigma(x_p)\}$ of $\{x_1, \ldots, x_n\}$, we deduce that $|\operatorname{Orb}(L)| \leq \binom{n}{p}$.

On the other hand, $\sigma \in \operatorname{Stab}(L)$ is equivalent to $\{\sigma(x_1), \ldots, \sigma(x_p)\} = \{x_1, \ldots, x_n\}$, which is further equivalent to $(X - \sigma(x_1)) \ldots (X - \sigma(x_p)) = (X - x_1) \ldots (X - x_p)$. Thus $\{\sigma(S_1), \ldots, \sigma(S_p)\} = \{S_1, \ldots, S_p\}$ so $\sigma(F) = F$, hence $\sigma \in G_1$. We conclude that $|G/\operatorname{Stab}(L)| = |G/G_1| = [F:K]$. Thus, it follows that $[F:K] \leq \binom{n}{p}$.

Solution by Marian Tetiva, Gheorghe Roşca Codreanu National College, Bârlad, Romania. For $1 \leq i_1 < i_2 < \cdots < i_p \leq n$ let us denote by

$$S_1(x_{i_1}, x_{i_2}, \dots, x_{i_p}), S_2(x_{i_1}, x_{i_2}, \dots, x_{i_p}), \dots, S_p(x_{i_1}, x_{i_2}, \dots, x_{i_p})$$

the symmetric fundamental polynomials in $x_{i_1}, x_{i_2}, \ldots, x_{i_p}$.

Thus, $S_1 = S_1(x_1, x_2, ..., x_p)$ and so on.

Let x be an arbitrary element from $K(S_1, S_2, \ldots, S_n)$; that is,

$$x = \sum a_{j_1 j_2 \dots j_p} S_1^{j_1} S_2^{j_2} \cdots S_p^{j_p},$$

for some $a_{j_1j_2...j_p} \in K$ (the sum being extended over a finite number of indices $0 \leq j_1 \leq n, 0 \leq j_2 \leq n, ..., 0 \leq j_p \leq n$). Denote this x by $x = x_{12...p}$ and define, similarly, $x_{i_1i_2...i_p}$ to be

$$\sum a_{j_1 j_2 \dots j_p} \left(S_1 \left(x_{i_1}, x_{i_2}, \dots, x_{i_p} \right) \right)^{j_1} \left(S_2 \left(x_{i_1}, x_{i_2}, \dots, x_{i_p} \right) \right)^{j_2} \cdots \cdots \left(S_p \left(x_{i_1}, x_{i_2}, \dots, x_{i_p} \right) \right)^{j_p}$$

with the same coefficients as in x, for all $1 \le i_1 < i_2 < \cdots < i_p \le n$. Now observe that the polynomial

$$f = \prod_{1 \le i_1 < i_2 < \dots < i_p \le n} (X - x_{i_1 i_2 \dots i_p})$$

has degree $\binom{n}{p}$, has coefficients in K (due to the fundamental theorem of symmetric polynomials), and has $x = x_{12...p}$ as a root.

Thus, any element of the algebraic extension $K \subseteq K(S_1, S_2, \dots, S_p)$ has degree over K at most $\binom{n}{p}$. This means that the degree of this extension is at most $\binom{n}{p}$, too, finishing the proof. \Box **Remark.** It seems that the condition about x_1, x_2, \ldots, x_n to be distinct is unnecessary.

306. Let $x_1, x_2, \ldots, x_n, y_1, y_2, \ldots, y_n$ be positive real numbers such that $x_1 + x_2 + \cdots + x_n = 1$, and denote by \mathcal{A} the set of pairs (i, j) such that $y_1 + \cdots + y_j \leq x_1 + \cdots + x_i$. Prove that one has

$$\sum_{(i,j)\in\mathcal{A}} \frac{x_{i+1}y_j}{1 + (x_1 + \dots + x_{i+1})(y_1 + \dots + y_j)} \le \frac{\pi^2}{24}$$

Prove that the constant $\frac{\pi^2}{24}$ is the best satisfying the property. Proposed by Radu Gologan, Simion Stoilow Institute of Mathematics

Proposed by Radu Gologan, Simion Stoilow Institute of Mathematics of the Romanian Academy, Bucharest, Romania.

Solution by the author. In the Cartesian plane consider the triangle T with vertices O(0,0), A(1,0), B(1,1). For $(i,j) \in \mathcal{A}$ put $a_i = x_1 + \cdots + x_i$, $b_j = y_1 + \cdots + y_j$, and denote by R_{ij} the rectangle with vertices of coordinates

$$(a_i, b_{j-1}), (a_i, b_j), (a_{i+1}, b_{j-1}), (a_{i+1}, b_j).$$

It is clear that $\bigcup_{\mathcal{A}} R_{ij}$ is contained in T and the rectangles have disjoint.

interiors.

Then the lower *Darboux* sum for $f(x, y) = \frac{1}{1 + xy}$ on *T* is less than

$$\sum_{R_{ij}} \min_{R_{ij}} f(x, y) x_{i+1} y_j \ge \sum_{(i,j) \in \mathcal{A}} \frac{x_{i+1} y_j}{1 + (x_1 + \dots + x_{i+1})(y_1 + \dots + y_j)}$$

by the fact that f is decreasing in each variable.

In conclusion, the last sum is less than the double integral

$$c = \iint_{T} \frac{1}{1 + xy} \mathrm{d}x \mathrm{d}y.$$

By Fubini calculation $c = \int_{0}^{1} \frac{\ln(1+x^2)}{x} dx$, so the existence part of the

result is proven. To prove that $c = \frac{\pi^2}{24}$, use the power series for $\ln(1+t)$, and the uniform convergence by the Abel theorem, on [0, 1]. Thus

$$\int_{0}^{1} \frac{\ln(1+x^2)}{x} dx = \int_{0}^{1} \sum_{n=1}^{\infty} (-1)^{n+1} x^n dx = \frac{1}{2} \sum_{n=1}^{\infty} (-1)^{n+1} \frac{1}{n^2} = \frac{\pi^2}{24},$$

concluding the proof.

To prove the fact that the found constant is the best, it suffices to consider in the inequality, for every n, the numbers $x_i = \frac{1}{n} = y_j$, and remark that the sum in the inequality is a Riemannian sum for the double integral, so for n going to infinity the sum tends to $\frac{\pi^2}{24}$.

307. Let $f : \mathbb{R} \to \mathbb{R}$ be twice differentiable with f'' continuous and $\lim_{x \to \pm \infty} f(x) = \infty$ such that f(x) > f''(x) for all $x \in \mathbb{R}$. Show that f(x) > 0 for any real number x.

Proposed by Beniamin Bogoşel, student West University of Timişoara, Timişoara, Romania.

Solution by Richard Stevens, Columbus State University, GA, USA. This solution follows also in large lines the solution of the author. Suppose that $f(a) \leq 0$ for some real number a. From the given conditions we see that there is an interval containing a for which f has a minimum value that does not occur at either end point of the interval. Denoting this minimum point as (b, f(b)), it follows that $f(b) \leq 0$, f'(b) = 0 and $f''(b) = \lim_{x \to b} \frac{f'(x)}{x - b} < 0$. Thus, for some $\varepsilon > 0$, f'(x) < 0 for $b < x < b + \varepsilon$ and f'(x) > 0 for $b - \varepsilon < x < b$. This indicates that (b, f(b)) is a relative maximum point and that f is not constant on an interval containing b. Therefore, f(x) > 0 for all x.

308. Let $M_n(\mathbb{Q})$ be the ring of square matrices of size n and $X \in M_n(\mathbb{Q})$. Define the adjugate (classical adjoint) of X, denoted $\operatorname{adj}(X)$, as follows. The (i, j)-minor M_{ij} of X is the determinant of the $(n-1) \times (n-1)$ matrix obtained by deleting row i and column j of X, and the (i, j) – cofactor of X is $C_{ij} = (-1)^{i+j} M_{ij}$. The adjugate of X is the transpose of the 'cofactor matrix' C_{ij} of X. Consider $A, B \in M_n(\mathbb{Q})$ such that

$$(\mathrm{adj}(A))^3 - (\mathrm{adj}(B))^3 = 2((\mathrm{adj}(A)) - (\mathrm{adj}(B))) \neq O_n.$$

Show that

$$\operatorname{rank}(AB) \in \{\operatorname{rank}(A), \operatorname{rank}(B)\}$$

Proposed by Flavian Georgescu, student University of Bucharest, Bucharest, Romania.

Solution by the author. We have to prove that either A or B is invertible. We shall argue by contradiction, so assume neither is invertible.

Firstly let us note that for non-invertible matrix $Y \in M_n(\mathbb{Q})$, we have rank $(\operatorname{adj}(Y)) \in \{0,1\}$. This is true since if rank $(Y) \leq n-2$, then $\operatorname{adj}(Y) = O_n$, otherwise if rank(Y) = n-1, using *Sylvester*'s inequality we have

 $\operatorname{rank}(Y \cdot \operatorname{adj}(Y)) + n \ge \operatorname{rank}(Y) + \operatorname{rank}(\operatorname{adj}(Y)),$

and since $Y \cdot \operatorname{adj}(Y) = O_n$, it follows that $\operatorname{rank}(\operatorname{adj}(Y)) \leq 1$. From above we get for A and B that $\operatorname{rank}(\operatorname{adj}(A))$, $\operatorname{rank}(\operatorname{adj}(B)) \in \{0, 1\}$, so we can deduce

that

$$(\operatorname{adj}(A))^2 = \operatorname{tr}(\operatorname{adj}(A)) \cdot \operatorname{adj}(A),$$

respectively

$$(\operatorname{adj}(B))^2 = \operatorname{tr}(\operatorname{adj}(B)) \cdot \operatorname{adj}(B).$$

For more ease, let us denote $\alpha = tr(adj(A))$ and $\beta = tr(adj(A))$. We can rewrite our hyphotesis as

$$\alpha^2 \cdot \operatorname{adj}(A) - \beta^2 \cdot \operatorname{adj}(B) = 2(\operatorname{adj}(A) - \operatorname{adj}(B)),$$

and passing to traces we get $\alpha^3 - \beta^3 = 2(\alpha - \beta)$, so

0

$$(\alpha - \beta)(\alpha^2 + \beta^2 + \alpha \cdot \beta - 2) = 0.$$

If $\alpha = \beta$, then $(\alpha^2 - 2)(\operatorname{adj}(A) - \operatorname{adj}(B)) = 0$, and since $\operatorname{adj}(A) \neq \operatorname{adj}(B)$, it would lead to $\alpha^2 - 2 = 0$, so $\sqrt{2} \in \mathbb{Q}$, a contradiction. If $\alpha^2 + \beta^2 + \alpha \cdot \beta = 2$, we can rewrite it as $(\alpha + 2\beta)^2 + 3\alpha^2 = 8$. We can reduce to the following equation in integers, $a^2 + 3b^2 = 8c^2$. This implies $3 \mid a^2 + c^2$ whence, since -1 is not a quadratic residue modulo 3, it follow that $3\mid a$ and $3\mid c$, thus $3\mid b$. Next we proceed by infinite descent, to obtain that the only solution is a = b = c = 0, a contradiction.

Thus our assumption is false, so one of the matrices A or B is invertible.

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